

Weak Dynamic Programming for Generalized State Constraints

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Abstract

We provide a dynamic programming principle for stochastic optimal control problems with expectation constraints. A weak formulation, using test functions and a probabilistic relaxation of the constraint, avoids restrictions related to a measurable selection but still implies the Hamilton-Jacobi-Bellman equation in the viscosity sense. We treat open state constraints as a special case of expectation constraints and prove a comparison theorem to obtain the equation for closed state constraints.

Keywords weak dynamic programming, state constraint, expectation constraint, Hamilton-Jacobi-Bellman equation, viscosity solution, comparison theorem

AMS 2000 Subject Classifications 93E20, 49L20, 49L25, 35K55

1 Introduction

We study the problem of stochastic optimal control under state constraints. In the most classical case, this is the problem of maximizing an expected reward, subject to the constraint that the controlled state process has to remain in a given subset \mathcal{O} of the state space. There is a rich literature on the associated partial differential equations (PDEs), going back to [13, 14, 9, 16, 17] in the first order case and [12, 11, 10] in the second order case. The connection between the control problem and the equation is given by the dynamic programming principle. However, in the stochastic case, it is frequent practice to make this connection only formally, and in fact, we are not aware of a generally applicable, rigorous technique in the literature. Of course, there are specific situations where it is indeed possible to avoid proving the state-constrained dynamic programming principle; in particular, penalization arguments can be useful to reduce to the unconstrained case (e.g., [11]). We refer to [6, 18] for further background.

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Generally speaking, it is difficult to prove the dynamic programming principle when the regularity of the value function is not known a priori, due to certain measurable selection problems. It was observed in [3] that, in the unconstrained case, these difficulties can be avoided by a weak formulation of the dynamic programming principle where the value function is replaced by a test function. This formulation, which is tailored to the derivation of the PDE in the sense of viscosity solutions, avoids the measurable selection and uses only a simple covering argument. It turns out that the latter does not extend directly to the case with state constraints. Essentially, the reason is that if ν is some admissible control for the initial condition $x \in \mathcal{O}$ —i.e., the controlled state process $X_{t,x}^\nu$ started at x remains in \mathcal{O} —then ν may fail to have this property for a nearby initial condition $x' \in \mathcal{O}$.

However, if \mathcal{O} is open and mild continuity assumptions are satisfied, then $X_{t,x'}^\nu$ will violate the state constraint with at most small probability when x' is close to x . This observation leads us to consider optimization problems with constraints in probability, and more generally expectation constraints of the form $E[g(X_{t,x}^\nu(T))] \leq m$ for given $m \in \mathbb{R}$. We shall see that, following the idea of [2], such problems are amenable to dynamic programming if the constraint level m is formulated dynamically via an auxiliary family of martingales. A key insight of the present paper is that relaxing the level m by a small constant allows to prove a weak dynamic programming principle for general expectation constraint problems (Theorem 2.4), while the PDE can be derived despite the relaxation. We shall then obtain the dynamic programming principle for the classical state constraint problem (Theorem 3.1) by passing to a limit $m \downarrow 0$, with a suitable choice of g and certain regularity assumptions. Of course, expectation constraints are of independent interest and use; e.g., in Mathematical Finance, where one considers the problem of maximizing expected terminal wealth $E[X_{0,x}^\nu(T)]$ under the loss constraint $E[(X_{0,x}^\nu(T) - x)^-]^p \leq m$, for some given $m, p > 0$.

We exemplify the use of these results in the setting of controlled diffusions and show how the PDEs for expectation constraints and state constraints can be derived (Theorems 4.2 and 4.6). For the latter case, we introduce an appropriate continuity condition at the boundary, under which we prove a comparison theorem. While the above concerned an open set \mathcal{O} and does not apply directly to the closed domain $\overline{\mathcal{O}}$, we show via the comparison result that the value function for $\overline{\mathcal{O}}$ coincides with the one for \mathcal{O} , under certain conditions.

The remainder of the paper is organized as follows. In Section 2 we introduce an abstract setup for dynamic programming under expectation constraints and prove the corresponding relaxed weak dynamic programming principle. In Section 3 we deduce the dynamic programming principle for the state constraint \mathcal{O} . We specialize to the case of controlled diffusions in Section 4, where we study the Hamilton-Jacobi-Bellman PDEs for expectation and state constraints. Appendix A provides the comparison theorem.

Throughout this paper, (in)equalities between random variables are in the almost sure sense and relations between processes are up to evanescence, unless otherwise stated.

2 Dynamic Programming Principle for Expectation Constraints

We fix a time horizon $T \in (0, \infty)$ and a probability space (Ω, \mathcal{F}, P) equipped with a filtration $\mathbb{F} = (\mathcal{F}_t)_{t \in [0, T]}$. For each $t \in [0, T]$, we are given a set \mathcal{U}_t whose elements are seen as controls at time t . Given a separable metric space S , the *state space*, we denote by $\mathbf{S} := [0, T] \times S$ the time-augmented state space. For each $(t, x) \in \mathbf{S}$ and $\nu \in \mathcal{U}_t$, we are given a càdlàg adapted process $X_{t,x}^\nu = \{X_{t,x}^\nu(s), s \in [t, T]\}$ with values in S , the controlled *state process*. Finally, we are given two measurable functions $f, g : S \rightarrow \mathbb{R}$. We assume that

$$E[|f(X_{t,x}^\nu(T))|] < \infty \quad \text{and} \quad E[|g(X_{t,x}^\nu(T))|] < \infty \quad \text{for all } \nu \in \mathcal{U}_t, \quad (2.1)$$

so that the *reward and constraint functions*

$$F(t, x; \nu) := E[f(X_{t,x}^\nu(T))], \quad G(t, x; \nu) := E[g(X_{t,x}^\nu(T))]$$

are well defined. We also introduce the *value function*

$$V(t, x, m) := \sup_{\nu \in \mathcal{U}(t, x, m)} F(t, x; \nu), \quad (t, x, m) \in \hat{\mathbf{S}}, \quad (2.2)$$

where $\hat{\mathbf{S}} := \mathbf{S} \times \mathbb{R} \equiv [0, T] \times S \times \mathbb{R}$ and

$$\mathcal{U}(t, x, m) := \{\nu \in \mathcal{U}_t : G(t, x; \nu) \leq m\} \quad (2.3)$$

is the set of controls admissible at constraint level m . Here $\sup \emptyset := -\infty$.

The following observation is the heart of our approach: Given a control ν admissible at level m at the point x , relaxing the level m will make ν admissible in an entire neighborhood of x . This will be crucial for the covering arguments used below. We use the acronym u.s.c. (l.s.c.) to indicate upper (lower) semicontinuity.

Lemma 2.1. *Let $(t, x, m) \in \hat{\mathbf{S}}$, $\nu \in \mathcal{U}(t, x, m)$ and assume that the function $x' \mapsto G(t, x'; \nu)$ is u.s.c. at x . For each $\delta > 0$ there exists a neighborhood B of $x \in S$ such that $\nu \in \mathcal{U}(t, x', m' + \delta)$ for all $x' \in B$ and all $m' \geq m$.*

Proof. We have $G(t, x; \nu) \leq m$ by the definition (2.3) of $\mathcal{U}(t, x, m)$. In view of the upper semicontinuity, there exists a neighborhood B of x such that $G(t, x'; \nu) \leq m + \delta \leq m' + \delta$ for $x' \in B$; that is, $\nu \in \mathcal{U}(t, x', m' + \delta)$. \square

A control problem with an expectation constraint of the form (2.3) is not amenable to dynamic programming if we just consider a fixed level m . Extending an idea from [2, 1], we shall see that this changes if the constraint is formulated dynamically by using auxiliary martingales. To this end, we consider for each $t \in [0, T]$ a family $\mathcal{M}_{t,0}$ of càdlàg martingales $M = \{M(s), s \in [t, T]\}$ with initial value $M(t) = 0$. We also introduce

$$\mathcal{M}_{t,m} := \{m + M : M \in \mathcal{M}_{t,0}\}, \quad m \in \mathbb{R}.$$

We assume that, for all $(t, x) \in \mathbf{S}$ and $\nu \in \mathcal{U}_t$,

$$\text{there exists } M_t^\nu[x] \in \mathcal{M}_{t,m} \text{ such that } M_t^\nu[x](T) = g(X_{t,x}^\nu(T)), \quad (2.4)$$

where, necessarily, $m = E[g(X_{t,x}^\nu(T))]$. In particular, given $\nu \in \mathcal{U}_t$, the set

$$\mathcal{M}_{t,m,x}^+(\nu) := \{M \in \mathcal{M}_{t,m} : M(T) \geq g(X_{t,x}^\nu(T))\}$$

is nonempty for any $m \geq E[g(X_{t,x}^\nu(T))]$. More precisely, we have the following characterization.

Lemma 2.2. *Let $(t, x) \in \mathbf{S}$ and $m \in \mathbb{R}$. Then*

$$\mathcal{U}(t, x, m) = \{\nu \in \mathcal{U}_t : \mathcal{M}_{t,m,x}^+(\nu) \neq \emptyset\}.$$

Proof. Let $\nu \in \mathcal{U}_t$. If there exists some $M \in \mathcal{M}_{t,m,x}^+(\nu)$, then taking expectations yields $E[g(X_{t,x}^\nu(T))] \leq E[M(T)] = m$ and hence $\nu \in \mathcal{U}(t, x, m)$.

Conversely, let $\nu \in \mathcal{U}(t, x, m)$; i.e., we have $m' := E[g(X_{t,x}^\nu(T))] \leq m$. With $M_t^\nu[x]$ as in (2.4), $M := M_t^\nu[x] + m - m'$ is an element of $\mathcal{M}_{t,m,x}^+(\nu)$. \square

It will be useful in the following to consider for each $t \in [0, T]$ an auxiliary subfiltration $\mathbb{F}^t = (\mathcal{F}_s^t)_{s \in [0, T]}$ of \mathbb{F} such that $X_{t,x}^\nu$ and M are \mathbb{F}^t -adapted for all $x \in S$, $\nu \in \mathcal{U}_t$ and $M \in \mathcal{M}_{t,0}$. Moreover, we denote by \mathcal{T}^t the set of \mathbb{F}^t -stopping times with values in $[t, T]$.

The following assumption corresponds to one direction of the dynamic programming principle; cf. Theorem 2.4(i) below.

Assumption A. For all $(t, x, m) \in \hat{\mathbf{S}}$, $\nu \in \mathcal{U}(t, x, m)$, $M \in \mathcal{M}_{t,m,x}^+(\nu)$, $\tau \in \mathcal{T}^t$ and P -a.e. $\omega \in \Omega$, there exists $\nu_\omega \in \mathcal{U}(\tau(\omega), X_{t,x}^\nu(\tau)(\omega), M(\tau)(\omega))$ such that

$$E[f(X_{t,x}^\nu(T)) | \mathcal{F}_\tau](\omega) \leq F(\tau(\omega), X_{t,x}^\nu(\tau)(\omega); \nu_\omega). \quad (2.5)$$

Next, we state two variants of the assumptions for the converse direction of the dynamic programming principle; we shall comment on the differences in Remark 2.5 below. In the first variant, the intermediate time is deterministic; this will be enough to cover stopping times with countably many values in Theorem 2.4(ii) below. We recall the notation $M_t^\nu[x]$ from (2.4).

Assumption B. Let $(t, x) \in \mathbf{S}$, $\nu \in \mathcal{U}_t$, $s \in [t, T]$, $\bar{\nu} \in \mathcal{U}_s$ and $\Gamma \in \mathcal{F}_s^t$.

(B1) There exists a control $\tilde{\nu} \in \mathcal{U}_t$ such that

$$X_{t,x}^{\tilde{\nu}}(\cdot) = X_{t,x}^\nu(\cdot) \quad \text{on } [t, T] \times (\Omega \setminus \Gamma); \quad (2.6)$$

$$X_{t,x}^{\tilde{\nu}}(\cdot) = X_{s, X_{t,x}^\nu(s)}^{\bar{\nu}}(\cdot) \quad \text{on } [s, T] \times \Gamma; \quad (2.7)$$

$$E[f(X_{t,x}^{\tilde{\nu}}(T)) | \mathcal{F}_s] \geq F(s, X_{t,x}^\nu(s); \bar{\nu}) \quad \text{on } \Gamma. \quad (2.8)$$

The control $\tilde{\nu}$ is denoted by $\nu \otimes_{(s, \Gamma)} \bar{\nu}$ and called a *concatenation* of ν and $\bar{\nu}$ on (s, Γ) .

(B2) Let $M \in \mathcal{M}_{t,0}$. There exists a process $\bar{M} = \{\bar{M}(r), r \in [s, T]\}$ such that

$$\bar{M}(\cdot)(\omega) = (M_s^{\bar{\nu}}[X_{t,x}^\nu(s)(\omega)](\cdot))(\omega) \quad \text{on } [s, T] \quad P\text{-a.s.}$$

and

$$M \mathbf{1}_{[t,s)} + \mathbf{1}_{[s,T]} \left(M \mathbf{1}_{\Omega \setminus \Gamma} + [\bar{M} - \bar{M}(s) + M(s)] \mathbf{1}_\Gamma \right) \in \mathcal{M}_{t,0}.$$

(B3) Let $m \in \mathbb{R}$ and $M \in \mathcal{M}_{t,m,x}^+(\nu)$. For P -a.e. $\omega \in \Omega$, there exist a control $\nu_\omega \in \mathcal{U}(s, X_{t,x}^\nu(s)(\omega), M(s)(\omega))$.

In the second variant, the intermediate time is a stopping time and we have an additional assumption about the structure of the sets \mathcal{U}_s . This variant corresponds to Theorem 2.4(ii') below.

Assumption B'. Let $(t, x) \in \mathbf{S}$, $\nu \in \mathcal{U}_t$, $\tau \in \mathcal{T}^t$, $\Gamma \in \mathcal{F}_\tau^t$ and $\bar{\nu} \in \mathcal{U}_{\|\tau\|_{L^\infty}}$.

(B0') $\mathcal{U}_s \supseteq \mathcal{U}_{s'}$ for all $0 \leq s \leq s' \leq T$.

(B1') There exists a control $\tilde{\nu} \in \mathcal{U}_t$, denoted by $\nu \otimes_{(\tau, \Gamma)} \bar{\nu}$, such that

$$X_{t,x}^{\tilde{\nu}}(\cdot) = X_{t,x}^\nu(\cdot) \quad \text{on } [t, T] \times (\Omega \setminus \Gamma);$$

$$X_{t,x}^{\tilde{\nu}}(\cdot) = X_{\tau, X_{t,x}^\nu(\tau)}^{\bar{\nu}}(\cdot) \quad \text{on } [\tau, T] \times \Gamma; \quad (2.9)$$

$$E[f(X_{t,x}^{\tilde{\nu}}(T)) | \mathcal{F}_\tau] \geq F(\tau, X_{t,x}^\nu(\tau); \bar{\nu}) \quad \text{on } \Gamma. \quad (2.10)$$

(B2') Let $M \in \mathcal{M}_{t,0}$. There exists a process $\bar{M} = \{\bar{M}(r), r \in [\tau, T]\}$ such that

$$\bar{M}(\cdot)(\omega) = (M_{\tau(\omega)}^{\bar{\nu}}[X_{t,x}^\nu(\tau)(\omega)](\cdot))(\omega) \quad \text{on } [\tau, T] \quad P\text{-a.s.}$$

and

$$M \mathbf{1}_{[t,\tau)} + \mathbf{1}_{[\tau,T]} \left(M \mathbf{1}_{\Omega \setminus \Gamma} + [\bar{M} - \bar{M}(\tau) + M(\tau)] \mathbf{1}_\Gamma \right) \in \mathcal{M}_{t,0}.$$

(B3') Let $m \in \mathbb{R}$ and $M \in \mathcal{M}_{t,m,x}^+(\nu)$. For P -a.e. $\omega \in \Omega$, there exist a control $\nu_\omega \in \mathcal{U}(\tau(\omega), X_{t,x}^\nu(\tau)(\omega), M(\tau)(\omega))$.

Remark 2.3. (i) Assumption B' implies Assumption B. Moreover, Assumption A implies (B3').

(ii) Let $\mathbf{D} := \{(t, x, m) \in \hat{\mathbf{S}} : \mathcal{U}(t, x, m) \neq \emptyset\}$ denote the natural domain of our optimization problem. Then (B3') can be stated as follows: for any $\nu \in \mathcal{U}(t, x, m)$ and $M \in \mathcal{M}_{t, m, x}^+(\nu)$,

$$(\tau, X_{t,x}^\nu(\tau), M(\tau)) \in \mathbf{D} \quad P\text{-a.s.} \quad \text{for all } \tau \in \mathcal{T}^t. \quad (2.11)$$

Under (B3), the same holds if τ takes countably many values. We remark that the invariance property (2.11) corresponds to one direction in the geometric dynamic programming of [15].

(iii) We note that (2.9) (and similarly (2.7)) states in particular that

$$(r, \omega) \mapsto (X_{\tau, X_{t,x}^\nu(\tau)}^{\bar{\nu}}(r))(\omega) := (X_{\tau(\omega), X_{t,x}^\nu(\tau)(\omega)}^{\bar{\nu}}(r))(\omega)$$

is a well-defined adapted process (up to evanescence) on $[t, T] \times \Gamma$. Of course, this is an implicit measurability condition on $(t, x) \mapsto X_{t,x}^{\bar{\nu}}$.

(iv) For an illustration of (B1'), let us assume that the controls are predictable processes. In this case, one can often take

$$\nu \otimes_{(\tau, \Gamma)} \bar{\nu} := \nu \mathbf{1}_{[0, \tau]} + \mathbf{1}_{(\tau, T]} (\bar{\nu} \mathbf{1}_\Gamma + \nu \mathbf{1}_{\Omega \setminus \Gamma}) \quad (2.12)$$

and the condition that $\nu \otimes_{(\tau, \Gamma)} \bar{\nu} \in \mathcal{U}_t$ is called *stability under concatenation*. The idea is that we use the control ν up to time τ ; after time τ , we continue using ν in the event $\Omega \setminus \Gamma$ while we switch to the control $\bar{\nu}$ in the event Γ . Of course, one can omit the set $\Gamma \in \mathcal{F}_\tau^t$ by observing that

$$\nu \otimes_{(\tau, \Gamma)} \bar{\nu} = \nu \mathbf{1}_{[0, \tau']} + \bar{\nu} \mathbf{1}_{(\tau', T]}$$

for the \mathbb{F}^t -stopping time $\tau' := \tau \mathbf{1}_\Gamma + T \mathbf{1}_{\Omega \setminus \Gamma}$.

We can now state our weak dynamic programming principle for the stochastic control problem (2.2) with expectation constraint. The formulation is weak in two ways; namely, the value function is replaced by a test function and the constraint level m is relaxed by an arbitrarily small constant $\delta > 0$ (cf. (2.15) below). The flexibility of choosing the set \mathcal{D} appearing below, will be used in Section 3. We recall the set \mathbf{D} introduced in Remark 2.3(ii).

Theorem 2.4. Let $(t, x, m) \in \hat{\mathbf{S}}$, $\nu \in \mathcal{U}(t, x, m)$ and $M \in \mathcal{M}_{t, m, x}^+(\nu)$. Let $\tau \in \mathcal{T}^t$ and let $\mathcal{D} \subseteq \hat{\mathbf{S}}$ be a set such that $(\tau, X_{t,x}^\nu(\tau), M(\tau)) \in \mathcal{D}$ holds P -a.s.

(i) Let Assumption A hold true and let $\varphi : \hat{\mathbf{S}} \rightarrow [-\infty, \infty]$ be a measurable function such that $V \leq \varphi$ on \mathcal{D} . Then $E[\varphi(\tau, X_{t,x}^\nu(\tau), M(\tau))]^- < \infty$ and

$$F(t, x; \nu) \leq E[\varphi(\tau, X_{t,x}^\nu(\tau), M(\tau))]. \quad (2.13)$$

(ii) Let $\delta > 0$, let Assumption B hold true and assume that τ takes countably many values $(t_i)_{i \geq 1}$. Let $\varphi : \hat{\mathbf{S}} \rightarrow [-\infty, \infty)$ be a measurable function such that $V \geq \varphi$ on \mathcal{D} . Assume that for any fixed $\bar{\nu} \in \mathcal{U}_{t_i}$,

$$\left. \begin{array}{ll} (x', m') \mapsto \varphi(t_i, x', m') & \text{is u.s.c.} \\ x' \mapsto F(t_i, x'; \bar{\nu}) & \text{is l.s.c.} \\ x' \mapsto G(t_i, x'; \bar{\nu}) & \text{is u.s.c.} \end{array} \right\} \text{on } \mathcal{D}^i \quad (2.14)$$

for all $i \geq 1$, where $\mathcal{D}^i := \{(x', m') : (t_i, x', m') \in \mathcal{D}\} \subseteq S \times \mathbb{R}$. Then

$$V(t, x, m + \delta) \geq E[\varphi(\tau, X_{t,x}^\nu(\tau), M(\tau))]. \quad (2.15)$$

(ii') Let $\delta > 0$ and let Assumption B' hold true. Let $\varphi : \hat{\mathbf{S}} \rightarrow [-\infty, \infty)$ be a measurable function such that $V \geq \varphi$ on \mathcal{D} . Assume that $\mathcal{D} \cap \mathbf{D} \subseteq \hat{\mathbf{S}}$ is σ -compact and that for any fixed $\bar{\nu} \in \mathcal{U}_{t_0}$, $t_0 \in [t, T]$,

$$\left. \begin{array}{ll} (t', x', m') \mapsto \varphi(t', x', m') & \text{is u.s.c.} \\ (t', x') \mapsto F(t', x'; \bar{\nu}) & \text{is l.s.c.} \\ (t', x') \mapsto G(t', x'; \bar{\nu}) & \text{is u.s.c.} \end{array} \right\} \text{on } \mathcal{D} \cap \{t' \leq t_0\}. \quad (2.16)$$

Then (2.15) holds true.

The following convention is used on the right hand side of (2.15): if Y is any random variable,

$$\text{we set } E[Y] := -\infty \quad \text{if } E[Y^+] = E[Y^-] = \infty. \quad (2.17)$$

We note that Assumption (B0') ensures that the expressions $F(t', x'; \bar{\nu})$ and $G(t', x'; \bar{\nu})$ in (2.16) are well defined for $t' \leq t_0$.

Remark 2.5. The difference between parts (ii) and (ii') of the theorem stems from the fact that in the proof of (ii') we consider $[0, T] \times S \times \mathbb{R}$ as the state space while for (ii) it suffices to consider $S \times \mathbb{R}$ and hence no assumptions on the time variable are necessary. Regarding applications, (ii) is obviously the better choice if stopping times with countably many values (and in particular deterministic times) are sufficient.

There is a number of cases where the extension to a general stopping time τ can be accomplished a posteriori by approximation, in particular if one has a priori estimates for the value function so that one can restrict to test functions with specific growth properties. Assume for illustration that f is bounded from above, then so is V and one will be interested only in test functions φ which are bounded from above; moreover, it will typically not hurt to assume that φ is u.s.c. (or even continuous) in all variables. Now let (τ_n) be a sequence of stopping time taking finitely many (e.g., dyadic) values such that $\tau_n \downarrow \tau$ P -a.s. Applying (2.15) to each τ^n and using Fatou's

lemma as well as the right-continuity of the paths of $X_{t,x}^\nu$ and M , we then find that (2.15) also holds for the general stopping time τ .

On the other hand, it is not always possible to pass to the limit as above and then (ii') is necessary to treat general stopping times. The compactness assumption should be reasonable *provided that* S is σ -compact; e.g., $S = \mathbb{R}^d$. Then, the typical way to apply (ii') is to take $(t, x, m) \in \text{int } \mathbf{D}$ and let \mathcal{D} be a open or closed neighborhood of (t, x, m) such that $\mathcal{D} \subseteq \mathbf{D}$.

Proof of Theorem 2.4. (i) Fix $(t, x, m) \in \hat{\mathbf{S}}$. With ν_ω as in (2.5), the definition (2.2) of V and $V \leq \varphi$ on \mathcal{D} imply that

$$\begin{aligned} E[f(X_{t,x}^\nu(T))|\mathcal{F}_\tau](\omega) &\leq F(\tau(\omega), X_{t,x}^\nu(\tau)(\omega); \nu_\omega) \\ &\leq V(\tau(\omega), X_{t,x}^\nu(\tau)(\omega), M(\tau)(\omega)) \\ &\leq \varphi(\tau(\omega), X_{t,x}^\nu(\tau)(\omega), M(\tau)(\omega)) \quad P\text{-a.s.} \end{aligned}$$

After noting that the left hand side is integrable by (2.1), the result follows by taking expectations.

(ii) Fix $(t, x, m) \in \hat{\mathbf{S}}$ and let $\varepsilon > 0$.

1. *Countable Cover.* Like \mathcal{D} , the set $\mathcal{D} \cap \mathbf{D}$ has the property that

$$(\tau, X_{t,x}^\nu(\tau), M(\tau)) \in \mathcal{D} \cap \mathbf{D} \quad P\text{-a.s.}; \quad (2.18)$$

this follows from (2.11) since τ takes countably many values. Replacing \mathcal{D} by $\mathcal{D} \cap \mathbf{D}$ if necessary, we may therefore assume that $\mathcal{D} \subseteq \mathbf{D}$. Using also that $V \geq \varphi$ on \mathcal{D} , the definition (2.2) of V shows that there exists for each $(s, y, n) \in \mathcal{D}$ some $\nu^{s,y,n} \in \mathcal{U}(s, y, n)$ satisfying

$$F(s, y; \nu^{s,y,n}) \geq \varphi(s, y, n) - \varepsilon. \quad (2.19)$$

Fix one of the points t_i in time. For each $(y, n) \in \mathcal{D}^i$, the semicontinuity assumptions (2.14) and Lemma 2.1 imply that there exists a neighborhood $(y, n) \in B^i(y, n) \subseteq \hat{S} := S \times \mathbb{R}$ (of size depending on $y, n, i, \varepsilon, \delta$) such that

$$\left. \begin{array}{l} \nu^{t_i, y, n} \in \mathcal{U}(t_i, y', n' + \delta) \\ \varphi(t_i, y', n') \leq \varphi(t_i, y, n) + \varepsilon \\ F(t_i, y'; \nu^{t_i, y, n}) \geq F(t_i, y; \nu^{t_i, y, n}) - \varepsilon \end{array} \right\} \text{ for all } (y', n') \in B^i(y, n) \cap \mathcal{D}^i. \quad (2.20)$$

Here the first inequality may read $-\infty \leq -\infty$. We note that $\mathcal{D}^i \subseteq \hat{S}$ is metric separable for the subspace topology relative to the product topology on \hat{S} . Therefore, since the family $\{B^i(y, n) \cap \mathcal{D}^i : (y, n) \in \hat{S}\}$ forms an open cover of \mathcal{D}^i , there exists a sequence $(y_j, n_j)_{j \geq 1}$ in \hat{S} such that $\{B^i(y_j, n_j) \cap \mathcal{D}^i\}_{j \geq 1}$ is a countable subcover of \mathcal{D}^i . We set $\nu_j^i := \nu^{t_i, y_j, n_j}$ and $B_j^i := B^i(y_j, n_j)$, so that

$$\mathcal{D}^i \subseteq \bigcup_{j \geq 1} B_j^i. \quad (2.21)$$

We can now define, for i still being fixed, a measurable partition $(A_j^i)_{j \geq 1}$ of $\cup_{j \geq 1} B_j^i$ by

$$A_1^i := B_1^i, \quad A_{j+1}^i := B_{j+1}^i \setminus (B_1^i \cup \dots \cup B_j^i), \quad j \geq 1.$$

Since $A_j^i \subseteq B_j^i$, the inequalities (2.19) and (2.20) yield that

$$F(t_i, y'; \nu_j^i) \geq \varphi(t_i, y', n') - 3\varepsilon \quad \text{for all } (y', n') \in A_j^i \cap \mathcal{D}^i. \quad (2.22)$$

2. *Concatenation.* Fix an integer $k \geq 1$; we now focus on $(t_i)_{1 \leq i \leq k}$. We may assume that $t_1 < t_2 < \dots < t_k$, by eliminating and relabeling some of the t_i . We define the \mathcal{F}_τ^t -measurable sets

$$\Gamma_j^i := \{\tau = t_i \text{ and } (X_{t,x}^\nu(t_i), M(t_i)) \in A_j^i\} \in \mathcal{F}_{t_i} \quad \text{and} \quad \Gamma(k) := \bigcup_{1 \leq i, j \leq k} \Gamma_j^i.$$

Since the t_i are distinct and $A_j^i \cap A_{j'}^i = \emptyset$ for $j \neq j'$, we have $\Gamma_j^i \cap \Gamma_{j'}^i = \emptyset$ for $(i, j) \neq (i', j')$. We can then consider the successive concatenation

$$\nu(k) := \nu \otimes_{(t_1, \Gamma_1^1)} \nu_1^1 \otimes_{(t_1, \Gamma_2^1)} \nu_2^1 \cdots \otimes_{(t_1, \Gamma_k^1)} \nu_k^1 \otimes_{(t_2, \Gamma_1^2)} \nu_1^2 \cdots \otimes_{(t_k, \Gamma_k^k)} \nu_k^k,$$

which is to be read from the left with $\nu^a \otimes \nu^b \otimes \nu^c := (\nu^a \otimes \nu^b) \otimes \nu^c$. It follows from an iterated application of (B1) that $\nu(k)$ is well defined and in particular $\nu(k) \in \mathcal{U}_t$. (To understand the meaning of $\nu(k)$, it may be useful to note that in the example considered in (2.12), we would have

$$\nu(k) = \nu \mathbf{1}_{[0, \tau]} + \mathbf{1}_{(\tau, T]} \left(\nu \mathbf{1}_{\Omega \setminus \Gamma(k)} + \sum_{1 \leq i, j \leq k} \nu_j^i \mathbf{1}_{\Gamma_j^i} \right);$$

i.e., at time τ , we switch to ν_j^i on Γ_j^i .) We note that (2.6) implies that $X_{t,x}^{\nu(k)} = X_{t,x}^\nu$ on $\Omega \setminus \Gamma(k) \in \mathcal{F}_\tau$ and hence

$$E[f(X_{t,x}^{\nu(k)}(T)) | \mathcal{F}_\tau] = E[f(X_{t,x}^\nu(T)) | \mathcal{F}_\tau] \quad \text{on } \Omega \setminus \Gamma(k). \quad (2.23)$$

Moreover, the fact that $\tau = t_i$ on Γ_j^i , repeated application of (2.6), and (2.8) show that

$$E[f(X_{t,x}^{\nu(k)}(T)) | \mathcal{F}_\tau] \geq F(\tau, X_{t,x}^\nu(\tau); \nu_j^i) \quad \text{on } \Gamma_j^i, \quad \text{for } 1 \leq i, j \leq k,$$

and we deduce via (2.22) that

$$\begin{aligned} E[f(X_{t,x}^{\nu(k)}(T)) | \mathcal{F}_\tau] \mathbf{1}_{\Gamma(k)} &\geq \sum_{1 \leq i, j \leq k} F(\tau, X_{t,x}^\nu(\tau); \nu_j^i) \mathbf{1}_{\Gamma_j^i} \\ &\geq \sum_{1 \leq i, j \leq k} (\varphi(\tau, X_{t,x}^\nu(\tau), M(\tau)) - 3\varepsilon) \mathbf{1}_{\Gamma_j^i} \\ &\geq \varphi(\tau, X_{t,x}^\nu(\tau), M(\tau)) \mathbf{1}_{\Gamma(k)} - 3\varepsilon. \end{aligned} \quad (2.24)$$

3. *Admissibility.* Next, we show that $\nu(k) \in \mathcal{U}(t, x, m + \delta)$. By (B2) there exists, for each pair $1 \leq i, j \leq k$, a process $M_j^i = \{M_j^i(u), u \in [t_i, T]\}$ such that

$$M_j^i(\cdot)(\omega) = \left(M_{t_i}^{\nu_j^i} [X_{t,x}^\nu(t_i)(\omega)](\cdot) \right) (\omega) \text{ on } \Gamma_j^i \quad (2.25)$$

and such that

$$\begin{aligned} M^{(k)} := & (M + \delta) \mathbf{1}_{[t,\tau]} \\ & + \mathbf{1}_{[\tau,T]} \left((M + \delta) \mathbf{1}_{\Omega \setminus \Gamma(k)} + \sum_{1 \leq i, j \leq k} [M_j^i - M_j^i(t_i) + M(t_i) + \delta] \mathbf{1}_{\Gamma_j^i} \right) \end{aligned}$$

is an element of $\mathcal{M}_{t,m+\delta}$. We note that $M^{(k)}(T) \geq M_j^i(T)$ on Γ_j^i since $M_j^i(t_i) \leq M(t_i) + \delta$ on Γ_j^i as a result of the first condition in (2.20). Hence, using (2.7) and (2.25), we have

$$g(X_{t,x}^{\nu(k)}(T)) = g(X_{t_i, X_{t,x}^\nu(t_i)}^{\nu_j^i}(T)) \leq M_j^i(T) \leq M^{(k)}(T) \quad \text{on } \Gamma_j^i. \quad (2.26)$$

This holds for all $1 \leq i, j \leq k$. On the other hand, using (2.6) and that $M \in \mathcal{M}_{t,m,x}^+(\nu)$, we have that

$$g(X_{t,x}^{\nu(k)}(T)) = g(X_{t,x}^\nu(T)) \leq M(T) \leq M(T) + \delta = M^{(k)}(T) \quad \text{on } \Omega \setminus \Gamma(k). \quad (2.27)$$

Combining (2.26) and (2.27), we have $g(X_{t,x}^{\nu(k)}(T)) \leq M^{(k)}(T)$ on Ω and so Lemma 2.2 yields that $\nu(k) \in \mathcal{U}(t, x, m + \delta)$.

4. *ε -Optimality.* We may assume that either the positive or the negative part of $\varphi(\tau, X_{t,x}^\nu(\tau), M(\tau))$ is integrable, as otherwise our claim (2.15) is trivial by (2.17). Using the definition (2.2) of V as well as (2.23) and (2.24), we have that

$$\begin{aligned} V(t, x, m + \delta) & \geq E[f(X_{t,x}^{\nu(k)}(T))] \\ & = E[E[f(X_{t,x}^{\nu(k)}(T)) | \mathcal{F}_\tau]] \\ & \geq E[\varphi(\tau, X_{t,x}^\nu(\tau), M(\tau)) \mathbf{1}_{\Gamma(k)}] - 3\varepsilon + E[f(X_{t,x}^\nu(T)) \mathbf{1}_{\Omega \setminus \Gamma(k)}] \end{aligned}$$

for every $k \geq 1$. Letting $k \rightarrow \infty$, we have $\Gamma(k) \uparrow \Omega$ P -a.s. by (2.18) and (2.21); therefore,

$$E[f(X_{t,x}^\nu(T)) \mathbf{1}_{\Omega \setminus \Gamma(k)}] \rightarrow 0$$

by dominated convergence and (2.1). Moreover, monotone convergence yields

$$E[\varphi(\tau, X_{t,x}^\nu(\tau), M(\tau)) \mathbf{1}_{\Gamma(k)}] \rightarrow E[\varphi(\tau, X_{t,x}^\nu(\tau), M(\tau))];$$

to see this, consider separately the cases where the positive or the negative part of $\varphi(\tau, X_{t,x}^\nu(\tau), M(\tau))$ are integrable. Hence we have shown that

$$V(t, x, m + \delta) \geq E[\varphi(\tau, X_{t,x}^\nu(\tau), M(\tau))] - 3\varepsilon.$$

As $\varepsilon > 0$ was arbitrary, this completes the proof of (ii).

(ii') Fix $(t, x, m) \in \hat{\mathbf{S}}$ and let $\varepsilon > 0$. In contrast to the proof of (ii), we shall cover a subset of $\hat{\mathbf{S}}$ rather than $S \times \mathbb{R}$. By (2.11), we may again assume that $\mathcal{D} \cap \mathbf{D} = \mathcal{D}$. Since $V \geq \varphi$ on \mathcal{D} , the definition (2.2) of V shows that there exists for each $(s, y, n) \in \mathcal{D}$ some $\nu^{s,y,n} \in \mathcal{U}(s, y, n)$ satisfying

$$F(s, y; \nu^{s,y,n}) \geq \varphi(s, y, n) - \varepsilon. \quad (2.28)$$

Given $(s, y, n) \in \mathcal{D}$, the semicontinuity assumptions (2.16) and a variant of Lemma 2.1 (including the time variable) imply that there exists a set $B(s, y, n) \subseteq \hat{\mathbf{S}}$ of the form

$$B(s, y, n) = ((s - r, s] \cap [0, T]) \times B_r(y, n),$$

where $r > 0$ and $B_r(y, n)$ is an open ball in $S \times \mathbb{R}_+$, such that

$$\left. \begin{array}{l} \nu^{s,y,n} \in \mathcal{U}(s', y', n' + \delta) \\ \varphi(s', y', n') \leq \varphi(s, y, n) + \varepsilon \\ F(s', y'; \nu^{s,y,n}) \geq F(s, y; \nu^{s,y,n}) - \varepsilon \end{array} \right\} \text{ for all } (s', y', n') \in B(s, y, n) \cap \mathcal{D}. \quad (2.29)$$

Note that we have exploited Assumption (B0'), which forces us to use the half-closed interval for the time variable. As a result, $B(s, y, n)$ is open for the product topology on $[0, T] \times S \times \mathbb{R}$ if S and \mathbb{R} are given the usual topology and $[0, T]$ is given the topology generated by the half-closed intervals (the upper limit topology). Let us denote the latter topological space by $[0, T]^*$.

Like the Sorgenfrey line, $[0, T]^*$ is a Lindelöf space. Let \mathcal{D}' be the canonical projection of $\mathcal{D} \subseteq \hat{\mathbf{S}} \equiv [0, T] \times S \times \mathbb{R}$ to $S \times \mathbb{R}$. Then \mathcal{D}' is again σ -compact. It is easy to see that the product of a Lindelöf space with a σ -compact space is again Lindelöf; in particular, $[0, T]^* \times \mathcal{D}'$ is Lindelöf. Moreover, since \mathcal{D} was σ -compact for the original topology on $\hat{\mathbf{S}}$ and the topology of $[0, T]^*$ is finer than the original one, we still have that $\mathcal{D} \subseteq [0, T]^* \times \mathcal{D}'$ is a countable union of closed subsets and therefore \mathcal{D} is Lindelöf also for the new topology.

By the Lindelöf property, the cover $\{B(s, y, n) \cap \mathcal{D} : (s, y, n) \in \hat{\mathbf{S}}\}$ of \mathcal{D} admits a countable subcover $\{B(s_j, y_j, n_j) \cap \mathcal{D}\}_{j \geq 1}$. We set $\nu_j := \nu^{s_j, y_j, n_j}$ and $B_j := B(s_j, y_j, n_j)$, then $\mathcal{D} \subseteq \bigcup_{j \geq 1} B_j$ and

$$A_1 := B_1, \quad A_{j+1} := B_{j+1} \setminus (B_1 \cup \dots \cup B_j), \quad j \geq 1$$

defines a measurable partition of $\bigcup_{j \geq 1} B_j$. Since $A_j \subseteq B_j$, the inequalities (2.28) and (2.29) yield that

$$F(s', y'; \nu_j) \geq \varphi(s', y', n') - 3\varepsilon \quad \text{for all } (s', y', n') \in A_j \cap \mathcal{D}.$$

Similarly as above, we first fix $k \geq 1$, define the \mathcal{F}_τ^t -measurable sets

$$\Gamma_j := \{(\tau, X_{t,x}^\nu(\tau), M(\tau)) \in A_j\} \quad \text{and} \quad \Gamma(k) := \bigcup_{1 \leq j \leq k} \Gamma_j$$

and set $\nu(k) := (\cdots ((\nu \otimes_{(\tau, \Gamma_1)} \nu_1) \otimes_{(\tau, \Gamma_2)} \nu_2) \cdots \otimes_{(\tau, \Gamma_k)} \nu_k)$. To check that

$$E[f(X_{t,x}^{\nu(k)}(T)) | \mathcal{F}_\tau] \geq F(\tau, X_{t,x}^\nu(\tau); \nu_j) \quad \text{on } \Gamma_j \quad \text{for } 1 \leq j \leq k,$$

we use that $\tau \leq s_j$ on Γ_j , so that we can apply (2.10) with the stopping time $\tilde{\tau} := \tau \wedge s_j$ satisfying $\|\tilde{\tau}\|_{L^\infty} \leq s_j$ and thus $\nu_j \in \mathcal{U}_{s_j} \subseteq \mathcal{U}_{\|\tilde{\tau}\|_{L^\infty}}$; c.f. (B0'). The rest of the proof is analogous to the above. \square

Remark 2.6. The assumption on σ -compactness in Theorem 2.4(ii') was used only to ensure the Lindelöf property of $\mathcal{D} \cap \mathbf{D}$ for the topology introduced in the proof. Therefore, any other assumption ensuring this will do as well.

Let us record a slight generalization of Theorem 2.4(ii),(ii') to the case of controls which are not necessarily admissible at the given point (t, x, m) . The intuition for this result is that the dynamic programming principle holds as before if we use such controls for a sufficiently short time (as formalized by condition (2.30) below) and then switch to admissible ones. More precisely, the proof also exploits the relaxation which is anyway present in (2.15). We use the notation of Theorem 2.4.

Corollary 2.7. *Let the assumptions of Theorem 2.4(ii) hold true except that $\nu \in \mathcal{U}_t$ and $M \in \mathcal{M}_{t,m}$ are not necessarily in $\mathcal{U}(t, x, m)$ and $\mathcal{M}_{t,m,x}^+(\nu)$, respectively. In addition, assume that*

$$(\tau, X_{t,x}^\nu(\tau), M(\tau)) \in \mathbf{D} \quad P\text{-a.s.} \quad (2.30)$$

Then the conclusion (2.15) of Theorem 2.4(ii) still holds true. Moreover, the same generalization holds true for Theorem 2.4(ii').

Proof. Let us inspect the proof of Theorem 2.4(ii). Using directly (2.30) rather than appealing to (2.11), the construction of the covering in Step 1 remains unchanged and the same is true for the concatenation in Step 2. In Step 3, we proceed as above up to and including (2.26). Note that (2.27) no longer holds as it used the assumption that $M \in \mathcal{M}_{t,m,x}^+(\nu)$. However, in view of (2.26) and $X_{t,x}^{\nu(k)}(T) = X_{t,x}^\nu(T)$ on $\Omega \setminus \Gamma(k)$, we still have

$$E[g(X_{t,x}^{\nu(k)}(T))] \leq E[M^{(k)}(T) \mathbf{1}_{\Gamma(k)}] + E[g(X_{t,x}^\nu(T)) \mathbf{1}_{\Omega \setminus \Gamma(k)}].$$

Since $g(X_{t,x}^\nu(T))$ is integrable by (2.1) and $\Gamma(k) \uparrow \Omega$, the latter expectation is bounded by δ for large k . Moreover, $\Gamma(k) \in \mathcal{F}_\tau$, the martingale property of $M^{(k)}$, and the fact that $M^{(k)} = M + \delta$ on $[0, \tau]$ yield that

$$E[M^{(k)}(T) \mathbf{1}_{\Gamma(k)}] = E[M^{(k)}(\tau) \mathbf{1}_{\Gamma(k)}] = E[(M(\tau) + \delta) \mathbf{1}_{\Gamma(k)}].$$

Since $E[M(\tau)] = m$, the right hand side is dominated by $m + 2\delta$ for large k . Together, we conclude that

$$E[g(X_{t,x}^{\nu(k)}(T))] \leq m + 3\delta;$$

i.e., $\nu(k) \in \mathcal{U}(t, x, m + 3\delta)$ for all large k . Step 4 of the previous proof then applies as before (recall that $f(X_{t,x}^\nu(T))$ is integrable by (2.1)), except that the changed admissibility of $\nu(k)$ now results in

$$V(t, x, m + 3\delta) \geq E[\varphi(\tau, X_{t,x}^\nu(\tau), M(\tau))].$$

However, since $\delta > 0$ was arbitrary, this is the same as (2.15). The argument to extend Theorem 2.4(ii') is analogous. \square

While we shall see that the relaxation $\delta > 0$ in (2.15) is harmless for the derivation of the Hamilton-Jacobi-Bellman equation, it is nevertheless interesting to know when the δ can be omitted; i.e., when V is right continuous in m . The following sufficient condition will be used when we consider state constraints.

Lemma 2.8. *Let $(t, x, m) \in \mathbf{D}$. For each $\delta > 0$ there is $\nu^\delta \in \mathcal{U}(t, x, m + \delta)$ such that*

$$F(t, x; \nu^\delta) \geq \delta^{-1} \wedge V(t, x, m + \delta) - \delta.$$

Let $\delta_0 > 0$ and assume that for all $0 < \delta \leq \delta_0$ there exists $\tilde{\nu}^\delta \in \mathcal{U}(t, x, m)$ such that

$$\lim_{\delta \downarrow 0} P\{X_{t,x}^{\nu^\delta}(T) \neq X_{t,x}^{\tilde{\nu}^\delta}(T)\} = 0$$

and such that the set

$$\{[f(X_{t,x}^{\nu^\delta}(T)) - f(X_{t,x}^{\tilde{\nu}^\delta}(T))]^+ : 0 < \delta \leq \delta_0\} \subseteq L^1(P) \quad (2.31)$$

is uniformly integrable. Then $m' \mapsto V(t, x, m')$ is right continuous at m .

Proof. Since $m' \mapsto \mathcal{U}(t, x, m')$ is increasing, we have that $(t, x, m) \in \mathbf{D}$ implies $\mathcal{U}(t, x, m + \delta) \neq \emptyset$ for all $\delta \geq 0$. Hence ν^δ exists; of course, the truncation at δ^{-1} is necessary only if $V(t, x, m + \delta) = \infty$. Moreover, the monotonicity of $m' \mapsto V(t, x, m')$ implies that the right limit $V(t, x, m+)$ exists and that $V(t, x, m+) \geq V(t, x, m)$; it remains to prove the opposite inequality. Let $0 < \delta \leq \delta_0$ and set $A^\delta := \{X_{t,x}^{\nu^\delta}(T) \neq X_{t,x}^{\tilde{\nu}^\delta}(T)\}$, then

$$\begin{aligned} V(t, x, m) &\geq F(t, x, \tilde{\nu}^\delta) \\ &= F(t, x, \nu^\delta) - E[1_{A^\delta}(f(X_{t,x}^{\nu^\delta}(T)) - f(X_{t,x}^{\tilde{\nu}^\delta}(T)))] \\ &\geq \delta^{-1} \wedge V(t, x, m + \delta) - \delta - E[1_{A^\delta}(f(X_{t,x}^{\nu^\delta}(T)) - f(X_{t,x}^{\tilde{\nu}^\delta}(T)))^+]. \end{aligned}$$

Letting $\delta \downarrow 0$, we deduce by (2.31) that $V(t, x, m) \geq V(t, x, m+)$. \square

Remark 2.9. The integrability assumption (2.31) is clearly satisfied if f is bounded. In the general case, it may be useful to consider the value function for a truncated function f in a first step.

Remark 2.10. Our results can be generalized to a setting with multiple constraints. Given $N \in \mathbb{N}$ and $m \in \mathbb{R}^N$, let

$$\mathcal{U}(t, x, m) := \{\nu \in \mathcal{U}_t : G^i(t, x; \nu) \leq m^i \text{ for } i = 1, \dots, N\},$$

where $G^i(t, x; \nu) := E[g^i(X_{t,x}^\nu(T))]$ for some measurable function g^i . In this case, $\mathcal{M}_{t,0}$ is defined as the set of càdlàg N -dimensional martingales $M = \{M(s), s \in [t, T]\}$ with initial value $M(t) = 0$, and

$$\mathcal{M}_{t,m} := \{m + M : M \in \mathcal{M}_{t,0}\}, \quad m \in \mathbb{R}^N.$$

This generalization, which has been considered in [4] within the framework of stochastic target problems with controlled loss, also allows to impose constraints at finitely many intermediate times $0 \leq T_1 \leq T_2 \leq \dots \leq T$. Indeed, we can increase the dimension of the state process and add the components $X_{t,x}^\nu(\cdot \wedge T_j)$.

3 Application to State Constraints

We consider an open set $\mathcal{O} \subseteq S := \mathbb{R}^d$ and study the stochastic control problem under the constraint that the state process has to stay in \mathcal{O} . Namely, we consider the value function

$$\bar{V}(t, x) := \sup_{\nu \in \bar{\mathcal{U}}(t, x)} F(t, x; \nu), \quad (t, x) \in \mathbf{S}, \quad (3.1)$$

where

$$\bar{\mathcal{U}}(t, x) := \{\nu \in \mathcal{U}_t : X_{t,x}^\nu(s) \in \mathcal{O} \text{ for all } s \in [t, T], P\text{-a.s.}\}.$$

In the following discussion we assume that, for all $(t, x) \in \mathbf{S}$ and $\nu \in \mathcal{U}_t$,

$$X_{t,x}^\nu \text{ has continuous paths;} \quad (3.2)$$

$$(t, x) \mapsto X_{t,x}^\nu(r) \text{ is continuous in probability, uniformly in } r; \quad (3.3)$$

$$\bar{\mathcal{U}}(t, x) \neq \emptyset \text{ for } (t, x) \in [0, T] \times \mathcal{O}. \quad (3.4)$$

Explicitly, the condition (3.3) means that $(t_n, x_n) \rightarrow (t, x)$ implies

$$\sup_{r \in [0, T]} d(X_{t_n, x_n}^\nu(r), X_{t,x}^\nu(r)) \rightarrow 0 \quad \text{in probability,}$$

where we set $X_{t,x}^\nu(r) := x$ for $r < t$, $d(\cdot, \cdot)$ denotes the Euclidean metric, and it is implicitly assumed that $\nu \in \mathcal{U}_{t_n}$ for all n . We shall augment the state process so that the state constraint becomes a special case of an expectation constraint. To this end, we introduce the distance function $d(x) := \inf\{d(x, x') : x' \in S \setminus \mathcal{O}\}$ for $x \in S$ and the auxiliary process

$$Y_{t,x,y}^\nu(s) := y \wedge \inf_{r \in [t, s]} d(X_{t,x}^\nu(r)), \quad s \in [t, T], \quad y \in [0, \infty). \quad (3.5)$$

By (3.2), each trajectory $\{X_{t,x}^\nu(r)(\omega), r \in [t, T]\} \subseteq S$ is compact; therefore, it has strictly positive distance to $S \setminus \mathcal{O}$ whenever it is contained in \mathcal{O} :

$$\{X_{t,x}^\nu(r)(\omega), r \in [t, T]\} \subseteq \mathcal{O} \quad \text{if and only if} \quad Y_{t,x,1}^\nu(T)(\omega) > 0.$$

We consider the augmented state process

$$\bar{X}_{t,x,y}^\nu(\cdot) := (X_{t,x}^\nu(\cdot), Y_{t,x,y}^\nu(\cdot))$$

on the state space $S \times [0, \infty)$, then

$$E[g(\bar{X}_{t,x,y}^\nu(T))] = P\{Y_{t,x,y}^\nu(T) \leq 0\} \quad \text{by setting} \quad g(x, y) := \mathbf{1}_{(-\infty, 0]}(y)$$

for $(x, y) \in S \times [0, \infty)$. Now the state constraint may be expressed as $E[g(\bar{X}_{t,x,1}^\nu(T))] \leq 0$ and therefore

$$\bar{\mathcal{U}}(t, x) = \mathcal{U}(t, x, 1, 0) \quad \text{and} \quad \bar{V}(t, x) = V(t, x, 1, 0);$$

of course, the value 1 may be replaced by any number $y > 0$. Here and in the sequel, we use the notation from the previous section applied to the controlled state process \bar{X} on $S \times [0, \infty)$; i.e., we tacitly replace the variable x by (x, y) to define the set $\bar{\mathcal{U}}(t, x, y, m)$ of admissible controls and the associated value function $\bar{V}(t, x, y, m)$.

One direction of the dynamic programming principle will be a consequence of the following counterpart of Assumption A.

Assumption \bar{A} . For all $(t, x) \in \mathbf{S}$, $\nu \in \bar{\mathcal{U}}(t, x)$, $\tau \in \mathcal{T}^t$ and P -a.e. $\omega \in \Omega$, there exists $\nu_\omega \in \bar{\mathcal{U}}(\tau(\omega), X_{t,x}^\nu(\tau)(\omega))$ such that

$$E[f(X_{t,x}^\nu(T)) | \mathcal{F}_\tau](\omega) \leq F(\tau(\omega), X_{t,x}^\nu(\tau)(\omega); \nu_\omega).$$

The more difficult direction of the dynamic programming principle will be inferred from Theorem 2.4 under a right-continuity condition; we shall exemplify in the subsequent section how to verify this condition.

Theorem 3.1. Consider $(t, x) \in \mathbf{S}$ and a family $\{\tau^\nu, \nu \in \bar{\mathcal{U}}(t, x)\} \subseteq \mathcal{T}^t$.

(i) Let Assumption \bar{A} hold true and let $\phi : \mathbf{S} \rightarrow [-\infty, \infty]$ be a measurable function such that $\bar{V} \leq \phi$. Then $E[\phi(\tau^\nu, X_{t,x}^\nu(\tau^\nu))^-] < \infty$ for all $\nu \in \bar{\mathcal{U}}(t, x)$ and

$$\bar{V}(t, x) \leq \sup_{\nu \in \bar{\mathcal{U}}(t, x)} E[\phi(\tau^\nu, X_{t,x}^\nu(\tau^\nu))].$$

(ii) Let Assumption B' hold true for the state process \bar{X} on $S \times [0, \infty)$ and let (3.2)–(3.4) hold true. Moreover, assume that

$$V(t, x, 1, 0) = V(t, x, 1, 0+) \tag{3.6}$$

and that $(s', x') \mapsto F(s', x'; \nu)$ is l.s.c. on $[0, t_0] \times \mathcal{O}$ for all $t_0 \in [t, T]$ and $\nu \in \mathcal{U}_{t_0}$. Then

$$\bar{V}(t, x) \geq \sup_{\nu \in \bar{\mathcal{U}}(t, x)} E[\phi(\tau^\nu, X_{t,x}^\nu(\tau^\nu))] \quad (3.7)$$

for any u.s.c. function $\phi : \mathbf{S} \rightarrow [-\infty, \infty)$ such that $\bar{V} \geq \phi$.

Proof. (i) We may assume that $\bar{\mathcal{U}}(t, x) \neq \emptyset$ as otherwise $\bar{V}(t, x) = -\infty$. As in the proof of (2.13), we obtain that $F(t, x; \nu) \leq E[\varphi(\tau, X_{t,x}^\nu(\tau))]$ for all $\nu \in \bar{\mathcal{U}}(t, x)$. The claim follows by taking supremum over ν .

(ii) Again, we may assume that $\bar{\mathcal{U}}(t, x) \neq \emptyset$ as otherwise the right hand side of (3.7) equals $-\infty$. We set

$$\mathcal{D} := [t, T] \times \mathcal{O} \times (0, \infty) \times \{0\}.$$

If $\nu \in \bar{\mathcal{U}}(t, x)$, then $g(\bar{X}_{t,x,1}^\nu(T)) = 0$ and hence the constant martingale $M := 0$ is contained in $\mathcal{M}_{t,0,(x,1)}^+(\nu)$. Moreover, $Y_{t,x,1}^\nu(s) > 0$ and hence $(s, \bar{X}_{t,x,1}^\nu(s), M(s)) \in \mathcal{D}$ for all $s \in [t, T]$, P -a.s. Furthermore, if we define

$$\varphi(t', x', y', m') := \begin{cases} \phi(t', x'), & (t', x', y', m') \in \mathcal{D} \\ -\infty, & \text{otherwise,} \end{cases}$$

then φ is u.s.c. on \mathcal{D} . To see that the third semicontinuity condition in (2.16) is also satisfied, note that for any $(t', x', y'), (t'', x'', y'') \in [0, t_0] \times S \times [0, \infty)$ and $\nu \in \mathcal{U}_{t_0}$,

$$\begin{aligned} & |Y_{t',x',y'}^\nu(T) - Y_{t'',x'',y''}^\nu(T)| \\ & \leq |y' - y''| + \left| \inf_{r \in [t', T]} d(X_{t',x'}^\nu(r)) - \inf_{r \in [t'', T]} d(X_{t'',x''}^\nu(r)) \right| \\ & \leq |y' - y''| + \sup_{r \in [0, T]} d(X_{t',x'}^\nu(r), X_{t'',x''}^\nu(r)). \end{aligned}$$

Hence (3.3) implies that $(t', x', y') \mapsto Y_{t',x',y'}^\nu(T)$ is continuous in probability. As $(-\infty, 0] \subset \mathbb{R}$ is closed, we conclude by the Portmanteau theorem that

$$(t', x', y') \mapsto P\{Y_{t',x',y'}^\nu(T) \in (\infty, 0]\} \equiv G(t', x', y'; \nu) \quad \text{is u.s.c.}$$

as required. Since any open subset of a Euclidean space is σ -compact and since (3.4) implies $\mathcal{D} \cap \mathbf{D} = \mathcal{D}$, we can use (3.6) and Theorem 2.4(ii') with $M = 0$ to obtain that

$$\begin{aligned} \bar{V}(t, x) &= V(t, x, 1, 0) \\ &= V(t, x, 1, 0+) \\ &\geq E[\varphi(\tau^\nu, X_{t,x}^\nu(\tau^\nu), Y_{t,x,1}^\nu(\tau^\nu), 0)] \\ &= E[\phi(\tau^\nu, X_{t,x}^\nu(\tau^\nu))]. \end{aligned}$$

As $\nu \in \bar{\mathcal{U}}(t, x)$ was arbitrary, the result follows. \square

Remark 3.2. Similarly as in Theorem 2.4(ii), there is also a version of Theorem 3.1(ii) for stopping times taking countably many values. In this case, Assumption B replaces Assumption B', all conditions on the time variable are superfluous, and one can consider a general separable metric space S .

4 Application to Controlled Diffusions

In this section, we show how the weak formulation of the dynamic programming principle applies in the context of controlled Brownian stochastic differential equations and how it allows to derive the Hamilton-Jacobi-Bellman PDEs for the value functions associated to optimal control problems with expectation or state constraints. As the main purpose of this section is to illustrate the use of Theorems 2.4 and 3.1, we shall choose a fairly simple setup allowing to explain the main points without too many distractions. Given the generality of those theorems, extensions such as singular control, mixed control/stopping problems, etc. do not present any particular difficulty.

4.1 Setup for Controlled Diffusions

From now on, we take $S = \mathbb{R}^d$ and let $\Omega = C([0, T]; \mathbb{R}^d)$ be the space of continuous paths, P the Wiener measure on Ω , and W the canonical process $W_t(\omega) = \omega_t$. Let $\mathbb{F} = (\mathcal{F}_t)_{t \in [0, T]}$ be the augmentation of the filtration generated by W ; without loss of generality, $\mathcal{F} = \mathcal{F}_T$. For $t \in [0, T]$, the auxiliary filtration $\mathbb{F}^t = (\mathcal{F}_s^t)_{s \in [0, T]}$ is chosen to be the augmentation of $\sigma(W_r - W_t, t \leq r \leq s)$; in particular, \mathbb{F}^t is independent of \mathcal{F}_t .

Consider a closed set $U \subseteq \mathbb{R}^d$ and let \mathcal{U} be the set of all U -valued predictable processes ν satisfying $E[\int_0^T |\nu_t|^2 dt] < \infty$. Then we set

$$\mathcal{U}_t = \{\nu \in \mathcal{U} : \nu \text{ is } \mathbb{F}^t\text{-predictable}\}.$$

This choice will be convenient to verify Assumption B'. We remark that the restriction to \mathbb{F}^t -predictable controls entails no loss of generality, in the sense that the alternative choice $\mathcal{U}_t = \mathcal{U}$ would result in the same value function. Indeed, this follows from a well known randomization argument (see, e.g., [3, Remark 5.2]).

Let \mathbb{M}^d denote the set of $d \times d$ matrices. Given two Lipschitz continuous functions

$$\mu : \mathbb{R}^d \times U \rightarrow \mathbb{R}^d, \quad \sigma : \mathbb{R}^d \times U \rightarrow \mathbb{M}^d$$

and $(t, x, \nu) \in [0, T] \times \mathbb{R}^d \times \mathcal{U}$, we define $X_{t,x}^\nu(\cdot)$ as the unique strong solution of the stochastic differential equation (SDE)

$$X(s) = x + \int_t^s \mu(X(r), \nu_r) dr + \int_t^s \sigma(X(r), \nu_r) dW_r, \quad t \leq s \leq T, \quad (4.1)$$

where we set $X_{t,x}^\nu(r) = x$ for $r \leq t$. As $X_{t,x}^\nu(T)$ is square integrable for any $\nu \in \mathcal{U}$, (2.1) is satisfied whenever

$$f \text{ and } g \text{ have quadratic growth,} \quad (4.2)$$

which we assume from now on. In addition, we also impose that

$$\begin{cases} f \text{ is l.s.c. and } f^- \text{ has subquadratic growth,} \\ g \text{ is u.s.c. and } g^+ \text{ has subquadratic growth,} \end{cases} \quad (4.3)$$

where $h : \mathbb{R}^d \rightarrow \mathbb{R}$ is said to have subquadratic growth if $h(x)/|x|^2 \rightarrow 0$ as $|x| \rightarrow \infty$. This will be used to obtain the semicontinuity properties (2.16).

Furthermore, we take $\mathcal{M}_{t,0}$ to be the family of all càdlàg martingales which start at 0 and are adapted to \mathbb{F}^t . By the independence of the increments of W , we see that $M \in \mathcal{M}_{t,0}$ is then also a martingale in the filtration \mathbb{F}^t and that $M_r = 0$ for $r \leq t$. For $\nu \in \mathcal{U}_t$, we have $X_{t,x}^\nu(T) \in L^2(\mathcal{F}_T^t, P)$ and hence (2.4) is satisfied with $M_t^\nu[x](\cdot) = E[X_{t,x}^\nu(T)|\mathcal{F}_\cdot^t]$. It will be useful to express the martingales as stochastic integrals. Let \mathcal{A}_t denote the set of \mathbb{R}^d -valued \mathbb{F}^t -predictable processes α such that $\int_0^T |\alpha_t|^2 dt < \infty$ P -a.s. and such that

$$M_{t,0}^\alpha(\cdot) := \int_t^\cdot \alpha_s^\top dW_s$$

is a martingale (\top denotes transposition). Then the Brownian representation theorem yields that

$$\mathcal{M}_{t,0} = \{M_{t,0}^\alpha : \alpha \in \mathcal{A}_t\}.$$

In the following, we also use the notation $M_{t,m}^\alpha := m + M_{t,0}^\alpha$.

Lemma 4.1. *In the above setup, Assumptions A, \bar{A} , B, B' are satisfied and F and G satisfy (2.16).*

Proof. Assumption (B0') is immediate from the definition of \mathcal{U}_t . We define the concatenation of controls by (2.12).

The validity of Assumptions A, \bar{A} and (B1') follows from the uniqueness and the flow property of (4.1); in particular, the control ν_ω in Assumption A can be defined by $\nu_\omega(\omega') := \nu(\omega \otimes_\tau \omega')$, $\omega' \in \Omega$, where the concatenated path $\omega \otimes_\tau \omega'$ is given by

$$(\omega \otimes_\tau \omega')_r := \omega_r \mathbf{1}_{[0,\tau(\omega)]}(r) + (\omega'_r - \omega'_\tau(\omega) + \omega_\tau(\omega)) \mathbf{1}_{(\tau(\omega),T]}(r).$$

While we refer to [3, Proposition 5.4] for a detailed proof, we emphasize that the choice of \mathcal{U}_s is crucial for the validity of (2.10): in the notation of Assumption B', (2.10) essentially requires that $\bar{\nu}$ be independent of \mathcal{F}_τ .

Let $t, \tau, \nu, \bar{\nu}$ be as in Assumption B', we show that (B2') holds. Let \bar{M} be a càdlàg version of

$$\bar{M}(r) := E[g(X_{t,x}^{\hat{\nu}}(T))|\mathcal{F}_r], \quad r \in [0, T], \quad \text{where } \hat{\nu} := \nu \mathbf{1}_{[0,\tau]} + \mathbf{1}_{(\tau,T]} \bar{\nu}.$$

By the same argument as in [3, Proposition 5.4], we deduce from the uniqueness and the flow property of (4.1) and the fact that $\bar{\nu}$ is independent of \mathcal{F}_τ that

$$E[g(X_{t,x}^{\hat{\nu}}(T))|\mathcal{F}_r] = E[g(X_{\tau,X_{t,x}^{\nu}(\tau)}^{\bar{\nu}}(T))|\mathcal{F}_r] = M_{\tau}^{\bar{\nu}}[X_{t,x}^{\nu}(\tau)](r) \quad \text{on } [\tau, T].$$

Hence $\bar{M} = M_{\tau}^{\bar{\nu}}[X_{t,x}^{\nu}(\tau)]$ on $[\tau, T]$. The last assertion of (B2') is clear by the definition of $\mathcal{M}_{t,0}$. As already mentioned in Remark 2.3, Assumption (B3') follows from Assumption A and Assumption B follows from Assumption B'.

Next, we check that F satisfies (2.16); i.e., that F is l.s.c. For fixed $\nu \in \mathcal{U}$, $(t, x) \mapsto X_{t,x}^{\nu}(T)$ is L^2 -continuous. Hence the semicontinuity from (4.3) and Fatou's lemma yield that $(t, x) \mapsto E[f(X_{t,x}^{\nu}(T))^+]$ is l.s.c. By the subquadratic growth from (4.3), we have that $\{f(X_{t,x}^{\nu}(T))^- : (t, x) \in B\}$ is uniformly integrable whenever $B \subset \mathbf{S}$ is bounded, hence the semicontinuity of f also yields that $(t, x) \mapsto E[f(X_{t,x}^{\nu}(T))^-]$ is u.s.c. As a result, F is l.s.c. The same arguments show that G also satisfies (2.16). \square

4.2 PDE for Expectation Constraints

In this section, we show how to deduce the Hamilton-Jacobi-Bellman PDE for the optimal control problem (2.2) with expectation constraint from the weak dynamic programming principle stated in Theorem 2.4. Given a suitably differentiable function $\varphi(t, x)$ on $[0, T] \times \mathbb{R}^d$, we shall denote by $\partial_t \varphi$ its derivative with respect to t and by $D\varphi$ and $D^2\varphi$ the Jacobian and the Hessian matrix with respect to x , respectively.

In the context of the setup introduced in the preceding Section 4.1, the Hamilton-Jacobi-Bellman operator is given by

$$H(x, p, Q) := \inf_{(u, a) \in U \times \mathbb{R}^d} (-L^{u,a}(x, p, Q)), \quad (x, p, Q) \in \mathbb{R}^d \times \mathbb{R}^{d+1} \times \mathbb{M}^{d+1},$$

where

$$L^{u,a}(x, p, Q) := \mu_{X,M}(x, u)^\top p + \frac{1}{2} \text{Tr}[\sigma_{X,M} \sigma_{X,M}^\top(x, u, a) Q], \quad (u, a) \in U \times \mathbb{R}^d$$

is the Dynkin operator with coefficients

$$\mu_{X,M}(x, u) := \begin{pmatrix} \mu(x, u) \\ 0 \end{pmatrix} \quad \text{and} \quad \sigma_{X,M}(x, u, a) := \begin{pmatrix} \sigma(x, u) \\ a^\top \end{pmatrix}.$$

Since the set $U \times \mathbb{R}^d$ is unbounded, H may be discontinuous and viscosity solution properties need to be stated in terms of the upper and lower semicontinuous envelopes of H ,

$$\begin{aligned} H^*(x, p, Q) &:= \limsup_{(x', p', Q') \rightarrow (x, p, Q)} H(x', p', Q'), \\ H_*(x, p, Q) &:= \liminf_{(x', p', Q') \rightarrow (x, p, Q)} H(x', p', Q'). \end{aligned}$$

The value function V defined in (2.2) may also be discontinuous and so we introduce

$$V^*(t, x, m) := \limsup_{\substack{(t', x', m') \rightarrow (t, x, m) \\ (t', x', m') \in \text{int } \mathbf{D}}} V(t', x', m'),$$

$$V_*(t, x, m) := \liminf_{\substack{(t', x', m') \rightarrow (t, x, m) \\ (t', x', m') \in \text{int } \mathbf{D}}} V(t', x', m').$$

Here $\text{int } \mathbf{D}$ denotes the parabolic interior; i.e, the interior of $\mathbf{D} \setminus \{t = T\}$ in $\hat{\mathbf{S}}$, where $\{t = T\} := \{(t, x, m) \in \hat{\mathbf{S}} : t = T\}$. Moreover, we shall denote by $\overline{\mathbf{D}}$ the closure of \mathbf{D} . The main result of this subsection is the following PDE.

Theorem 4.2. *Assume that V is locally bounded on $\text{int } \mathbf{D}$.*

(i) *The function V^* is a viscosity subsolution on $\overline{\mathbf{D}} \setminus \{t = T\}$ of*

$$-\partial_t \varphi + H_*(\cdot, D\varphi, D^2\varphi) \leq 0.$$

(ii) *The function V_* is a viscosity supersolution on $\text{int } \mathbf{D}$ of*

$$-\partial_t \varphi + H^*(\cdot, D\varphi, D^2\varphi) \geq 0.$$

We refer to [5] for the various equivalent definitions of viscosity super- and subsolutions. We merely mention that “subsolution on A ” means that the subsolution property is satisfied at points of A which are local maxima of $V^* - \varphi$ on A , where φ is a test function, and analogously for the supersolution.

We shall not discuss in this generality the boundary condition and the validity of a comparison principle. In the subsequent section, these will be studied in some detail for the case of state constraints. We also refer to [1] for the study of the boundary conditions in a similar framework.

Remark 4.3. We observe that the domain of the PDE in Theorem 4.2 is not given a priori; it is itself characterized by a control problem: if we define

$$v(t, x) := \inf_{\nu \in \mathcal{U}_t} E[g(X_{t,x}^\nu(T))], \quad (t, x) \in \mathbf{S}, \quad (4.4)$$

then

$$\text{int } \mathbf{D} = \{(t, x, m) \in \hat{\mathbf{S}} : m > v^*(t, x), t < T\},$$

where v^* is the upper semicontinuous envelope of v on $[0, T] \times \mathbb{R}^d$. In particular, $\text{int } \mathbf{D} \neq \emptyset$ since v is locally bounded from above. In fact, in the present setup, we also have

$$\text{int } \mathbf{D} = \{(t, x, m) \in \hat{\mathbf{S}} : m > v(t, x), t < T\}. \quad (4.5)$$

Indeed, a well known randomization argument (e.g., [3, Remark 5.2]) yields that $v(t, x) = \inf_{\nu \in \mathcal{U}} E[g(X_{t,x}^\nu(T))]$ for all $(t, x) \in \mathbf{S}$; i.e., the set \mathcal{U}_t in (4.4)

can be replaced by \mathcal{U} . Therefore, v inherits the upper semicontinuity of G (c.f. Lemma 4.1) and we have $v = v^*$. Using (4.5), we obtain that

$$\{(t, x, m) \in \hat{\mathbf{S}} : m \geq v(t, x)\} \subseteq \{(t, x, m) \in \hat{\mathbf{S}} : m \geq v_*(t, x)\} = \overline{\text{int } \mathbf{D}},$$

where v_* is the lower semicontinuous envelope of v on $[0, T] \times \mathbb{R}^d$, and hence

$$\mathbf{D} \subseteq \overline{\text{int } \mathbf{D}}. \quad (4.6)$$

If furthermore v is continuous and the infimum in (4.4) is attained for all (t, x) , then the converse inclusion is also satisfied. In applications, it may be desirable to have continuity of v so that $\text{int } \mathbf{D}$ and $\overline{\mathbf{D}}$ are described directly by v (rather than a semicontinuous envelope). To analyze the continuity of v , one can study the comparison principle for the Hamilton-Jacobi-Bellman equation associated with the control problem (4.4). We refer to [8, Sections 5 and 6] for the explicit computation of v in an example from Mathematical Finance.

The rest of this subsection is devoted to the proof of Theorem 4.2. We first state a version of Theorem 2.4 which is suitable to derive the PDE.

Lemma 4.4. (i) *Let B be an open neighborhood of a point $(t, x, m) \in \mathbf{D}$ such that $V(t, x, m) < \infty$ and let $\varphi : \overline{B} \rightarrow \mathbb{R}$ be a continuous function such that $V \leq \varphi$ on \overline{B} . For all $\varepsilon > 0$ there exist $(\nu, \alpha) \in \mathcal{U}_t \times \mathcal{A}_t$ such that*

$$V(t, x, m) \leq E[\varphi(\tau, X_{t,x}^\nu(\tau), M_{t,m}^\alpha(\tau))] + \varepsilon \quad \text{and} \quad M_{t,m}^\alpha(T) \geq g(X_{t,x}^\nu(T)),$$

where τ is the first exit time of $(s, X_{t,x}^\nu(s), M_{t,m}^\alpha(s))_{s \geq t}$ from B .

(ii) *Let B be an open neighborhood of a point $(t, x, m) \in \text{int } \mathbf{D}$ such that $\overline{B} \subseteq \mathbf{D}$ and let $(\nu, \alpha) \in \mathcal{U}_t \times \mathcal{A}_t$. For any continuous function $\varphi : \overline{B} \rightarrow \mathbb{R}$ satisfying $V \geq \varphi$ on \overline{B} and for any $\varepsilon > 0$,*

$$V(t, x, m + \varepsilon) \geq E[\varphi(\tau, X_{t,x}^\nu(\tau), M_{t,m}^\alpha(\tau))], \quad (4.7)$$

where τ is the first exit time of $(s, X_{t,x}^\nu(s), M_{t,m}^\alpha(s))_{s \geq t}$ from B .

Proof. In view of Lemmata 2.2 and 4.1, part (i) is immediate from Theorem 2.4(i).

For part (ii) we use the extension of Theorem 2.4(ii') as stated in Corollary 2.7 with $\mathcal{D} := \overline{B}$. Note that $(\tau, X_{t,x}^\nu(\tau), M_{t,x}^\alpha(\tau)) \in \overline{B} \subseteq \mathbf{D}$; in particular, $\mathcal{D} \cap \mathbf{D} = \mathcal{D}$ is closed and hence σ -compact. \square

We can now deduce the PDE for V from the dynamic programming principle in the form of Lemma 4.4. Although the arguments are the usual ones, we shall indicate the proof, in particular to show that the relaxation “ $m + \varepsilon$ ” in (4.7) does not affect the PDE.

Proof of Theorem 4.2. (i) We first prove the subsolution property. Let φ be a $C^{1,2}$ -function and let $(t_0, x_0, m_0) \in \overline{\mathbf{D}}$ be such that $t_0 \in (0, T)$ and (t_0, x_0, m_0) is a maximum point of $V^* - \varphi$ satisfying

$$(V^* - \varphi)(t_0, x_0, m_0) = 0. \quad (4.8)$$

Assume for contradiction that

$$(-\partial_t \varphi + H_*(\cdot, D\varphi, D^2\varphi))(t_0, x_0, m_0) > 0.$$

Since $\overline{\mathbf{D}} = \overline{\text{int } \mathbf{D}}$ by (4.6), there exists a bounded open neighborhood $B \subset \hat{\mathbf{S}}$ of (t_0, x_0, m_0) such that

$$-\partial_t \bar{\varphi} - L^{u,a}(\cdot, D\bar{\varphi}, D^2\bar{\varphi}) > 0 \quad \text{on } \overline{B} \cap \overline{\text{int } \mathbf{D}}, \quad \text{for all } (u, a) \in U \times \mathbb{R}^d, \quad (4.9)$$

where

$$\bar{\varphi}(s, y, n) := \varphi(t_0, x_0, m_0) + (|s - t_0|^2 + |y - x_0|^4 + |n - m_0|^4).$$

Moreover, we have

$$\eta := \min_{\partial B} (\bar{\varphi} - \varphi) > 0. \quad (4.10)$$

Given $\varepsilon > 0$, let $(t_\varepsilon, x_\varepsilon, m_\varepsilon) \in B \cap \text{int } \mathbf{D}$ be such that

$$V(t_\varepsilon, x_\varepsilon, m_\varepsilon) \geq V^*(t_0, x_0, m_0) - \varepsilon. \quad (4.11)$$

Consider arbitrary $(\nu, \alpha) \in \mathcal{U}_t \times \mathcal{A}_t$ such that $M_{t_\varepsilon, m_\varepsilon}^\alpha(T) \geq g(X_{t_\varepsilon, x_\varepsilon}^\nu(T))$ and let τ be the first exit time of $(s, X_{t_\varepsilon, x_\varepsilon}^\nu(s), M_{t_\varepsilon, m_\varepsilon}^\alpha(s))_{s \geq t}$ from B . We recall from Remark 2.3(ii) that $(s, X_{t_\varepsilon, x_\varepsilon}^\nu(s), M_{t_\varepsilon, m_\varepsilon}^\alpha(s))_{s \geq t}$ remains in \mathbf{D} on $[t, T]$, and hence also in $\overline{\text{int } \mathbf{D}}$ by (4.6). Now, it follows from Itô's formula and (4.9) that

$$\bar{\varphi}(t_\varepsilon, x_\varepsilon, m_\varepsilon) \geq E[\bar{\varphi}(\tau, X_{t_\varepsilon, x_\varepsilon}^\nu(\tau), M_{t_\varepsilon, m_\varepsilon}^\alpha(\tau))].$$

For $(t_\varepsilon, x_\varepsilon, m_\varepsilon)$ close enough to (t_0, x_0, m_0) , this implies that

$$\bar{\varphi}(t_0, x_0, m_0) \geq E[\bar{\varphi}(\tau, X_{t_\varepsilon, x_\varepsilon}^\nu(\tau), M_{t_\varepsilon, m_\varepsilon}^\alpha(\tau))] - o(1) \quad \text{as } \varepsilon \rightarrow 0,$$

which, by (4.8), (4.10) and (4.11), leads to

$$V(t_\varepsilon, x_\varepsilon, m_\varepsilon) \geq E[\varphi(\tau, X_{t_\varepsilon, x_\varepsilon}^\nu(\tau), M_{t_\varepsilon, m_\varepsilon}^\alpha(\tau))] + \eta - o(1).$$

This contradicts Lemma 4.4(i) for $\varepsilon > 0$ small enough.

(ii) We now prove the supersolution property. Let φ be a $C^{1,2}$ -function and let $(t_0, x_0, m_0) \in \text{int } \mathbf{D}$ be such that (t_0, x_0, m_0) is a minimum point of $V_* - \varphi$ satisfying

$$(V_* - \varphi)(t_0, x_0, m_0) = 0. \quad (4.12)$$

Assume for contradiction that

$$(-\partial_t \varphi + H^*(\cdot, D\varphi, D^2\varphi))(t_0, x_0, m_0) < 0.$$

Then there exist $(\hat{u}, \hat{a}) \in U \times \mathbb{R}^d$ and a bounded open neighborhood B of (t_0, x_0, m_0) such that $\overline{B} \subseteq \text{int } \mathbf{D}$ and

$$-\partial_t \bar{\varphi} - L^{\hat{u}, \hat{a}}(\cdot, D\bar{\varphi}, D^2\bar{\varphi}) < 0 \quad \text{on } B, \quad (4.13)$$

where

$$\bar{\varphi}(s, y, n) := \varphi(t_0, x_0, m_0) - (|s - t_0|^2 + |y - x_0|^4 + |n - m_0|^4).$$

Note that

$$\eta := \min_{\partial B} (\varphi - \bar{\varphi}) > 0. \quad (4.14)$$

Given $\varepsilon > 0$, let $(t_\varepsilon, x_\varepsilon, m_\varepsilon) \in B$ be such that

$$V(t_\varepsilon, x_\varepsilon, m_\varepsilon + \varepsilon) \leq V_*(t_0, x_0, m_0) + \varepsilon. \quad (4.15)$$

Viewing (\hat{u}, \hat{a}) as a constant control, it follows from Itô's formula and (4.13) that

$$\bar{\varphi}(t_\varepsilon, x_\varepsilon, m_\varepsilon) \leq E[\bar{\varphi}(\tau, X_{t_\varepsilon, x_\varepsilon}^{\hat{u}}(\tau), M_{t_\varepsilon, m_\varepsilon}^{\hat{a}}(\tau))],$$

where τ is the first exit time of $(s, X_{t_\varepsilon, x_\varepsilon}^{\hat{u}}(s), M_{t_\varepsilon, m_\varepsilon}^{\hat{a}}(s))_{s \geq t}$ from B . For $(t_\varepsilon, x_\varepsilon, m_\varepsilon)$ close enough to (t_0, x_0, m_0) , this implies that

$$\bar{\varphi}(t_0, x_0, m_0) \leq E[\bar{\varphi}(\tau, X_{t_\varepsilon, x_\varepsilon}^{\hat{u}}(\tau), M_{t_\varepsilon, m_\varepsilon}^{\hat{a}}(\tau))] + o(1) \quad \text{as } \varepsilon \rightarrow 0,$$

which, by (4.12), (4.14) and (4.15), leads to

$$V(t_\varepsilon, x_\varepsilon, m_\varepsilon + \varepsilon) \leq E[\varphi(\tau, X_{t_\varepsilon, x_\varepsilon}^{\hat{u}}(\tau), M_{t_\varepsilon, m_\varepsilon}^{\hat{a}}(\tau))] - \eta + o(1).$$

For $\varepsilon > 0$ small enough, this yields a contradiction to Lemma 4.4(ii). \square

4.3 PDE for State Constraints

In this section, we discuss the Hamilton-Jacobi-Bellman PDE for the state constraint problem (cf. Section 3) in the case where the state process is given by a controlled SDE as introduced in Section 4.1 and required to stay in an open set $\mathcal{O} \subseteq \mathbb{R}^d$. Note that in this setup, the continuity conditions (3.2) and (3.3) are satisfied.

We shall derive the PDE via Theorem 3.1. The basic idea to guarantee its condition (3.6) about right continuity in the constraint level runs as follows. Consider a control $\nu \in \mathcal{U}_t$ such that $X_{t,x}^\nu$ leaves \mathcal{O} with at most small probability $\delta \geq 0$. Then we shall construct a control ν^δ , satisfying the state constraint, by switching to some admissible control $\hat{\nu}$ shortly before $X_{t,x}^\nu$ exits \mathcal{O} . As a result, ν^δ coincides with ν on a set of large probability and therefore the reward is similar. Along the lines of Lemma 2.8 we shall then obtain the desired right continuity (cf. Lemma 4.7 below).

To make this work, we clearly need to have $\bar{\mathcal{U}}(t, x) \neq \emptyset$ for all (t, x) in $[0, T] \times \mathcal{O}$, which is anyway necessary for the value function \bar{V} from (3.1)

to be finite. However, we need a slightly stronger condition; namely, that we can switch to an admissible control in a measurable way. A particularly simple condition ensuring this, is the existence of an admissible feedback control:

Assumption C. There exists a Lipschitz continuous mapping $\hat{u} : \mathcal{O} \rightarrow U$ such that, for all $(t, x) \in [0, T] \times \mathcal{O}$, the solution $\hat{X}_{t,x}$ of

$$\hat{X}(s) = x + \int_t^s \mu(\hat{X}(r), \hat{u}(\hat{X}(r))) dr + \int_t^s \sigma(\hat{X}(r), \hat{u}(\hat{X}(r))) dW_r, \quad s \in [t, T] \quad (4.16)$$

satisfies $\hat{X}_{t,x}(s) \in \mathcal{O}$ for all $s \in [t, T]$, P -a.s.

If, e.g., $\mu(\cdot, u_0) = 0$ and $\sigma(\cdot, u_0) = 0$ for some $u_0 \in U$, then Assumption C is clearly satisfied for $\hat{u} \equiv u_0$. Or, under an additional smoothness condition, $\hat{X}_{t,x}$ will stay in \mathcal{O} if the Lipschitz function \hat{u} satisfies

$$|n\sigma|(\cdot, \hat{u}) = 0 \quad \text{and} \quad \left(n^\top \mu + \frac{1}{2} \text{Trace}[Dn \sigma \sigma^\top] \right) (\cdot, \hat{u}) > 0 \quad \text{on } \partial\mathcal{O},$$

where n denotes the inner normal to $\partial\mathcal{O}$; see also [11, Proposition 3.1] and [7, Lemma III.4.3].

The following is a simple condition guaranteeing the uniform integrability required in (2.31).

Assumption D. Either f is bounded or the coefficients $\mu(x, u)$ and $\sigma(x, u)$ in the SDE (4.1) have linear growth in x , uniformly in u .

This assumption holds in particular if the control domain U is bounded.

Remark 4.5. Assumption C implies that \bar{V} is locally bounded from below and Assumption D implies that \bar{V} is locally bounded from above.

Next, we introduce the notation for the PDE related to the value function \bar{V} from (3.1). The associated Hamilton-Jacobi-Bellman operator is given by

$$\bar{H}(x, p, Q) := \inf_{u \in U} (-\bar{L}^u(x, p, Q)), \quad (x, p, Q) \in \mathbb{R}^d \times \mathbb{R}^d \times \mathbb{M}^d, \quad (4.17)$$

where the Dynkin operator is defined by

$$\bar{L}^u(x, p, Q) := \mu(x, u)^\top p + \frac{1}{2} \text{Tr}[\sigma \sigma^\top(x, u)Q], \quad u \in U.$$

Similarly as above, we introduce the semicontinuous envelopes

$$\begin{aligned} \bar{H}^*(x, p, Q) &:= \limsup_{\substack{(x', p', Q') \rightarrow (x, p, Q) \\ (x', p', Q') \in \mathcal{O} \times \mathbb{R}^d \times \mathbb{M}^d}} \bar{H}(x', p', Q'), \\ \bar{H}_*(x, p, Q) &:= \liminf_{\substack{(x', p', Q') \rightarrow (x, p, Q) \\ (x', p', Q') \in \mathcal{O} \times \mathbb{R}^d \times \mathbb{M}^d}} \bar{H}(x', p', Q') \end{aligned}$$

and

$$\begin{aligned}\bar{V}^*(t, x) &:= \limsup_{\substack{(t', x') \rightarrow (t, x) \\ (t', x') \in [0, T) \times \mathcal{O}}} \bar{V}(t', x'), \\ \bar{V}_*(t, x) &:= \liminf_{\substack{(t', x') \rightarrow (t, x) \\ (t', x') \in [0, T) \times \mathcal{O}}} \bar{V}(t', x').\end{aligned}$$

We can now state the Hamilton-Jacobi-Bellman PDE; the boundary condition is discussed in Proposition 4.11 below.

Theorem 4.6. *Assume that \bar{V} is locally bounded on $[0, T) \times \mathcal{O}$.*

(i) *The function \bar{V}^* is a viscosity subsolution on $[0, T) \times \bar{\mathcal{O}}$ of*

$$-\partial_t \varphi + \bar{H}_*(\cdot, D\varphi, D^2\varphi) \leq 0. \quad (4.18)$$

(ii) *Under Assumptions C and D, the function \bar{V}_* is a viscosity supersolution on $[0, T) \times \mathcal{O}$ of*

$$-\partial_t \varphi + \bar{H}^*(\cdot, D\varphi, D^2\varphi) \geq 0.$$

The proof is given below, after some auxiliary results. We first verify the right-continuity condition (3.6) for the value function $V(t, x, y, m)$ introduced in Section 3.

Lemma 4.7. *Let Assumptions C and D hold true. Then $V(t, x, 1, 0+) = V(t, x, 1, 0)$ for all $(t, x) \in [0, T] \times \mathcal{O}$.*

Proof. For $\delta > 0$, let $\nu^\delta \in \mathcal{U}(t, x, \delta)$ be as in Lemma 2.8. Then the process $Y_{t,x,y}^{\nu^\delta}$ defined in (3.5) satisfies $Y_{t,x,1}^{\nu^\delta}(T) > 0$ outside of a set of measure at most δ . It follows that we can find $\varepsilon_\delta \in (0, 1 \wedge d(x))$ such that the set $A^\delta := \{Y_{t,x,1}^{\nu^\delta}(T) \leq \varepsilon_\delta\}$ satisfies $P[A^\delta] \leq 2\delta$. Let τ^δ denote the first time when $Y_{t,x,1}^{\nu^\delta}$ reaches the level ε_δ and set

$$\tilde{\nu}^\delta := \nu^\delta \mathbf{1}_{[t, \tau^\delta]} + \mathbf{1}_{(\tau^\delta, T]} \hat{u}(\hat{X}^\delta),$$

where \hat{X}^δ is the solution of (4.16) on $[\tau^\delta, T]$ with initial condition given by $\hat{X}^\delta(\tau^\delta) = X_{t,x}^{\nu^\delta}(\tau^\delta)$. Since the paths of $Y_{t,x,1}^{\nu^\delta}$ are nonincreasing, we have

$$\lim_{\delta \downarrow 0} P[X_{t,x}^{\tilde{\nu}^\delta}(T) \neq X_{t,x}^{\nu^\delta}(T)] \leq \lim_{\delta \downarrow 0} P[\tau^\delta \leq T] = \lim_{\delta \downarrow 0} P[A^\delta] = 0.$$

Next, we check that $\{f(X_{t,x}^\nu(T)), \nu \in \mathcal{U}\}$ is uniformly integrable. This is trivial if f is bounded. Otherwise, Assumption D yields that the coefficients $\mu(x, u)$ and $\sigma(x, u)$ have uniformly linear growth in x , and of course they are uniformly Lipschitz in x as they are jointly Lipschitz. Thus $\{X_{t,x}^\nu(T), \nu \in \mathcal{U}\}$ is bounded in L^p for any finite p and the uniform integrability follows from the quadratic growth assumption (4.2) on f . It remains to apply Lemma 2.8. \square

Lemma 4.8. *Let Assumption C hold true and let B be an open neighborhood of $(t, x) \in [0, T] \times \mathcal{O}$. For all $\nu \in \mathcal{U}_t$ there exists $\bar{\nu} \in \mathcal{U}_t$ such that*

$$\nu \mathbf{1}_{[t, \tau]} + \bar{\nu} \mathbf{1}_{(\tau, T]} \in \bar{\mathcal{U}}(t, x),$$

where τ is the first exit time of $(s, X_{t,x}^\nu(s))_{s \geq t}$ from B .

Proof. Let $\hat{X}_{\tau, X_{t,x}^\nu(\tau)}$ be the solution of (4.16) on $[\tau, T]$ with the (square integrable) initial condition $X_{t,x}^\nu(\tau)$ at time τ . Then the claim holds true for $\bar{\nu} := \nu \mathbf{1}_{[t, \tau]} + \mathbf{1}_{(\tau, T]} \hat{u}(\hat{X}_{\tau, X_{t,x}^\nu(\tau)})$. \square

We have the following counterpart of Lemma 4.4.

Lemma 4.9. (i) *Let $B \subseteq [0, T] \times \mathbb{R}^d$ be an open neighborhood of a point $(t, x) \in [0, T] \times \mathcal{O}$ such that $\bar{V}(t, x)$ is finite and let $\varphi : \bar{B} \rightarrow \mathbb{R}$ be a continuous function such that $\bar{V} \leq \varphi$ on \bar{B} . For all $\varepsilon > 0$ there exists $\nu \in \bar{\mathcal{U}}(t, x)$ such that*

$$\bar{V}(t, x) \leq E[\varphi(\tau, X_{t,x}^\nu(\tau))] + \varepsilon,$$

where τ is the first exit time of $(s, X_{t,x}^\nu(s))_{s \geq t}$ from B .

(ii) *Let Assumptions C and D hold true and let $B \subseteq [0, T] \times \mathcal{O}$ be an open neighborhood of (t, x) . For any $\nu \in \mathcal{U}_t$ and any continuous function $\varphi : \bar{B} \rightarrow \mathbb{R}$ satisfying $\bar{V} \geq \varphi$ on \bar{B} ,*

$$\bar{V}(t, x) \geq E[\varphi(\tau, X_{t,x}^\nu(\tau))],$$

where τ is the first exit time of $(s, X_{t,x}^\nu(s))_{s \geq t}$ from B .

Proof. In view of Lemma 4.1, part (i) is immediate from Theorem 3.1(i). Part (ii) follows from Theorem 3.1(ii) via Lemmata 4.7 and 4.8. \square

Proof of Theorem 4.6. The result follows from Lemma 4.9 by the arguments in the proof of Theorem 4.2. \square

4.3.1 Boundary Condition and Uniqueness

In this section, we discuss the boundary condition and the uniqueness for the PDE in Theorem 4.6. We shall work under a slightly stronger condition on our setup.

Assumption D'. The coefficients $\mu(x, u)$ and $\sigma(x, u)$ in the SDE (4.1) have linear growth in x , uniformly in u .

We also introduce the following regularity condition, which will be used to prove the comparison theorem.

Definition 4.10. Consider a set $\mathcal{O} \subseteq \mathbb{R}^d$ and a function $w : [0, T] \times \bar{\mathcal{O}} \rightarrow \mathbb{R}$. Then w is of class $\mathcal{R}(\mathcal{O})$ if the following hold for any $(t, x) \in [0, T] \times \partial\mathcal{O}$:

(i) There exist $r > 0$, an open neighborhood B of x in \mathbb{R}^d and a function $\ell : \mathbb{R}_+ \rightarrow \mathbb{R}^d$ such that

$$\liminf_{\varepsilon \rightarrow 0} \varepsilon^{-1} |\ell(\varepsilon)| < \infty \quad \text{and} \quad (4.19)$$

$$y + \ell(\varepsilon) + o(\varepsilon) \in \mathcal{O} \quad \text{for all } y \in \overline{\mathcal{O}} \cap B \text{ and } \varepsilon \in (0, r). \quad (4.20)$$

(ii) There exists a function $\lambda : \mathbb{R}_+ \rightarrow \mathbb{R}_+$ such that

$$\lim_{\varepsilon \rightarrow 0} \lambda(\varepsilon) = 0 \quad \text{and} \quad (4.21)$$

$$\lim_{\varepsilon \rightarrow 0} w(t + \lambda(\varepsilon), x + \ell(\varepsilon)) = w(t, x). \quad (4.22)$$

By (4.20) we mean that if $\psi : \mathbb{R}_+ \rightarrow \mathbb{R}^d$ is any function of class $o(\varepsilon)$, then there exists $r > 0$ such that $y + \ell(\varepsilon) + \psi(\varepsilon) \in \mathcal{O}$ for all $\varepsilon \in (0, r)$. Note that (i) is a condition on the boundary of $\partial\mathcal{O}$; it can be seen as a variant of the interior cone condition where the cone is replaced by a more general shape. Condition (ii) essentially states that w is continuous along at least one curve approaching the boundary point through this shape. In Proposition 4.12 below, we indicate a sufficient condition for \bar{V}_* to be of class $\mathcal{R}(\mathcal{O})$, which is stated directly in terms of the given primitives. We shall see in its proof that Definition 4.10 is well adapted to the problem at hand (see also Remark A.4 below). Before that, let us state the uniqueness result.

Proposition 4.11. *Let f be continuous and let Assumptions C and D' hold true. Then \bar{V} has quadratic growth and the boundary condition is attained in the sense that*

$$\bar{V}^*(T, \cdot) \leq f \quad \text{and} \quad \bar{V}_*(T, \cdot) \geq f \quad \text{on } \overline{\mathcal{O}}.$$

Assume in addition that \bar{V}_* is of class $\mathcal{R}(\mathcal{O})$. Then

- (i) \bar{V} is continuous on $[0, T] \times \mathcal{O}$ and admits a continuous extension to $[0, T] \times \overline{\mathcal{O}}$,
- (ii) \bar{V} is the unique (discontinuous) viscosity solution of the state constraint problem

$$-\partial_t \varphi + \bar{H}(\cdot, D\varphi, D^2\varphi) = 0, \quad \varphi(T, \cdot) = f$$

in the class of functions having polynomial growth and having a lower semicontinuous envelope of class $\mathcal{R}(\mathcal{O})$.

Proof. Recalling that f has quadratic growth (4.2), it follows from standard estimates for the SDE (4.1) under Assumptions C and D' that \bar{V} has quadratic growth and satisfies the boundary condition.

Assumption D' implies Assumption D and hence Theorem 4.6 yields that \bar{V}^* and \bar{V}_* are sub- and supersolutions, respectively. Moreover, Assumption D' implies that \bar{H} is continuous and satisfies Assumption E in

Appendix A; see Lemma A.3 below for the proof. If \bar{V}_* is of class $\mathcal{R}(\mathcal{O})$, the comparison principle (Theorem A.1 and the subsequent remark, cf. Appendix A) yields that $\bar{V}^* = \bar{V}_*$ on $[0, T] \times \bar{\mathcal{O}}$; in particular, (i) holds. Part (ii) also follows from the comparison result. \square

We conclude this section with a sufficient condition ensuring that \bar{V}_* is of class $\mathcal{R}(\mathcal{O})$; the idea is that the volatility should degenerate so that the state process can be pushed away from the boundary. We remark that conditions in a similar spirit exist in the previous literature (e.g., [11, 10]); cf. Remark A.4 below.

Proposition 4.12. *Assume that \bar{V}_* is finite-valued on $[0, T] \times \bar{\mathcal{O}}$ and that \mathcal{O} , μ and σ satisfy the following conditions:*

(i) *There exists a C^1 -function δ , defined on a neighborhood of $\bar{\mathcal{O}} \subseteq \mathbb{R}^d$, such that $D\delta$ is locally Lipschitz continuous and*

$$\delta > 0 \text{ on } \mathcal{O}, \quad \delta = 0 \text{ on } \partial\mathcal{O}, \quad \delta < 0 \text{ outside } \bar{\mathcal{O}}.$$

(ii) *There exists a locally Lipschitz continuous mapping $\check{u} : \mathbb{R}^d \rightarrow U$ such that for all $x \in \bar{\mathcal{O}}$ there exist an open neighborhood B of x and $\iota > 0$ satisfying*

$$\mu(z, \check{u}(z))^\top D\delta(y) \geq \iota \text{ and } \sigma(y, \check{u}(y)) = 0 \text{ for all } y \in B \cap \bar{\mathcal{O}} \text{ and } z \in B. \quad (4.23)$$

Then \bar{V}_ is of class $\mathcal{R}(\mathcal{O})$.*

Proof. Fix $(t, x) \in [0, T] \times \partial\mathcal{O}$ and let δ , ι and B be as above; we may assume that B is bounded. Consider $y \in \bar{\mathcal{O}} \cap B$. Since $x' \mapsto \mu(x', \check{u}(x'))$ is locally Lipschitz, the ordinary differential equation

$$\check{x}(s) = y + \int_0^s \mu(\check{x}(r), \check{u}(\check{x}(r))) dr$$

has a unique solution \check{x}_y on some interval $[0, T_y]$. Then (4.23) ensures that

$$\check{x}_y(s) \in \mathcal{O} \quad \text{for } s \in (0, \varepsilon], \text{ for } \varepsilon > 0 \text{ small enough, for all } y \in \bar{\mathcal{O}} \cap B. \quad (4.24)$$

For $\varepsilon \in [0, T_x]$, we set

$$\ell(\varepsilon) := \check{x}_x(\varepsilon) - x, \quad \lambda(\varepsilon) := \varepsilon.$$

Then (4.21) is clearly satisfied. Using the continuity of μ and \check{u} , we see that

$$\varepsilon^{-1} |\ell(\varepsilon)| \rightarrow \mu(x, \check{u}(x)),$$

which implies (4.19). Moreover, using that μ and \check{u} are locally Lipschitz, we find that

$$\check{x}_x(r) - \check{x}_y(r) = x - y + O(r).$$

Together with (4.23), (4.24) and the local Lipschitz continuity of μ, \check{u} and $D\delta$, this implies that for $y \in \overline{\mathcal{O}}$ sufficiently close to x and $\varepsilon > 0$ small enough,

$$\begin{aligned} & \delta(y + \ell(\varepsilon) + o(\varepsilon)) \\ &= \delta(y - x + \check{x}_x(\varepsilon)) + o(\varepsilon) \\ &= \delta(y) + \int_0^\varepsilon \mu(\check{x}_x(r), \check{u}(\check{x}_x(r)))^\top D\delta(y - x + \check{x}_x(r)) dr + o(\varepsilon) \\ &= \delta(y) + \int_0^\varepsilon \mu(x - y + \check{x}_y(r), \check{u}(x - y + \check{x}_y(r)))^\top D\delta(\check{x}_y(r)) dr + O(\varepsilon^2) + o(\varepsilon) \\ &\geq \varepsilon \iota + o(\varepsilon), \end{aligned}$$

which is strictly positive for $\varepsilon > 0$ small enough. This implies (4.20).

Consider $(s, y) \in [0, T) \times \mathcal{O}$ close to (t, x) . For $\varepsilon > 0$ small enough, we can find $\nu^\varepsilon \in \bar{\mathcal{U}}(s + \lambda(\varepsilon), \check{x}_y(\varepsilon))$ such that

$$E[f(X_{s+\lambda(\varepsilon), \check{x}_y(\varepsilon)}^{\nu^\varepsilon}(T))] \geq \bar{V}(s + \lambda(\varepsilon), \check{x}_y(\varepsilon)) - \varepsilon.$$

Recall the degeneracy condition in (4.23). Setting

$$\bar{\nu}^\varepsilon := \check{u}(\check{x}_y) \mathbf{1}_{[s, s+\lambda(\varepsilon)]} + \mathbf{1}_{(s+\lambda(\varepsilon), T]} \nu^\varepsilon,$$

we obtain that

$$\begin{aligned} \bar{V}(s, y) &\geq E[f(X_{s,y}^{\bar{\nu}^\varepsilon}(T))] \\ &= E[f(X_{s+\lambda(\varepsilon), \check{x}_y(\varepsilon)}^{\nu^\varepsilon}(T))] \\ &\geq \bar{V}(s + \lambda(\varepsilon), \check{x}_y(\varepsilon)) - \varepsilon. \end{aligned}$$

Recalling that $\check{x}_y(\varepsilon) \rightarrow \check{x}_x(\varepsilon)$ as $y \rightarrow x$, this leads to

$$\bar{V}_*(t, x) \geq \bar{V}_*(t + \lambda(\varepsilon), \check{x}_x(\varepsilon)) - \varepsilon = \bar{V}_*(t + \lambda(\varepsilon), x + \ell(\varepsilon)) - \varepsilon$$

for $\varepsilon > 0$ small enough, which implies

$$\begin{aligned} \bar{V}_*(t, x) &\geq \limsup_{\varepsilon \rightarrow 0} \bar{V}_*(t + \lambda(\varepsilon), x + \ell(\varepsilon)) \\ &\geq \liminf_{\varepsilon \rightarrow 0} \bar{V}_*(t + \lambda(\varepsilon), x + \ell(\varepsilon)) \\ &\geq \bar{V}_*(t, x). \end{aligned}$$

Hence $\lim_{\varepsilon \rightarrow 0} \bar{V}_*(t + \lambda(\varepsilon), x + \ell(\varepsilon)) = \bar{V}_*(t, x)$; i.e., (4.22) holds for \bar{V}_* . \square

4.3.2 On Closed State Constraints

Recall that the value function \bar{V} considered above corresponds to the constraint that the state processes remains in the open set \mathcal{O} . We can similarly consider the closed constraint; i.e.,

$$\bar{V}(t, x) := \sup \{E[f(X_{t,x}^\nu(T))]: \nu \in \mathcal{U}_t, X_{t,x}^\nu(s) \in \overline{\mathcal{O}} \text{ for all } s \in [t, T], P\text{-a.s.}\}$$

The arguments used above for \bar{V} do not apply to \bar{V} . Indeed, for the closed set, the constraint function G in Section 3 would not be u.s.c. and hence the derivation of Theorem 3.1 fails; note that the upper semicontinuity is essential for the covering argument in the proof of Theorem 2.4(ii),(ii'). Moreover, the switching argument in the proof of Lemma 4.8 cannot be imitated since, given that the state process $X_{t,x}^\nu$ hits the boundary $\partial\mathcal{O}$, it is not possible to know which trajectories of the state will actually exit $\bar{\mathcal{O}}$.

However, we shall see that, if a comparison principle holds, then the dynamic programming principle for the open constraint \mathcal{O} is enough to fully characterize the value function \bar{V} associated to $\bar{\mathcal{O}}$. More precisely, we shall apply the PDE for \bar{V} and its comparison principle to deduce that $\bar{V} = \bar{V}$ under certain conditions. Of course, the basic observation that this equality holds under suitable conditions is not new; see, e.g., [10]. We shall use the following assumption.

Assumption C'. Assumption C holds with \hat{u} defined on $\bar{\mathcal{O}}$.

Corollary 4.13. *Let f be continuous, let Assumptions C' and D' hold true and assume that \bar{V}_* is of class $\mathcal{R}(\mathcal{O})$. Then $\bar{V} = \bar{V}_*$ on $[0, T] \times \bar{\mathcal{O}}$.*

We recall that \bar{V} admits a (unique) continuous extension to $[0, T] \times \bar{\mathcal{O}}$ under the stated conditions, so the assertion makes sense.

Proof. The easier direction of the dynamic programming principle for \bar{V} can be obtained as above, so the result of Lemma 4.9(i) still holds. It then follows by the same arguments as in the proof of Theorem 4.6 and Proposition 4.11 that the upper semicontinuous envelope \bar{V}^* is a viscosity subsolution of (4.18) satisfying $\bar{V}^*(T, \cdot) \leq f$ on $\bar{\mathcal{O}}$. As in the proof of Proposition 4.11, we can apply the comparison principle of Theorem A.1 to deduce that $\bar{V}^* \leq \bar{V}_*$ on $[0, T] \times \bar{\mathcal{O}}$. On the other hand, we clearly have $\bar{V} \leq \bar{V}$ on $[0, T] \times \mathcal{O}$ by the definitions of these value functions. Therefore, we have

$$\bar{V}_* \leq \bar{V}_* \leq \bar{V}^* \leq \bar{V}_* = \bar{V}^*,$$

where \bar{V}_* denotes the lower semicontinuous envelope of \bar{V} and the last equality is due to Proposition 4.11. It follows that all these functions coincide. \square

A Comparison for State Constraint Problems

In this appendix we provide, by adapting the usual techniques, a fairly general comparison theorem for state constraint problems which is suitable for the applications in Proposition 4.11 and Corollary 4.13.

In the following, \mathcal{H} denotes a continuous mapping from $\mathbb{R}^d \times \mathbb{R}^d \times \mathbb{M}^d$ to \mathbb{R} which is nonincreasing in its third variable, \mathcal{O} is a given open subset of \mathbb{R}^d , and $\rho > 0$ is a fixed constant. We consider the equation

$$\rho\varphi - \partial_t\varphi + \mathcal{H}(\cdot, D\varphi, D^2\varphi) = 0 \tag{A.1}$$

and the following condition on \mathcal{H} .

Assumption E. There exists $\alpha > 0$ such that

$$\begin{aligned} \liminf_{\eta \downarrow 0} (\mathcal{H}(y, q, Y^\eta) - \mathcal{H}(x, p, X^\eta)) \\ \leq \alpha \left(|x - y| (1 + |q| + n^2 |x - y|) + (1 + |x|) |p - q| + (1 + |x|^2) |Q| \right) \end{aligned}$$

for all $(x, y) \in \overline{\mathcal{O}}$ with $|x - y| \leq 1$ and for all $(p, q, Q) \in \mathbb{R}^d \times \mathbb{R}^d \times \mathbb{M}^{2d}$, $(X^\eta, Y^\eta)_{\eta > 0} \subset \mathbb{M}^d \times \mathbb{M}^d$ and $n \geq 1$ such that

$$\begin{pmatrix} X^\eta & 0 \\ 0 & -Y^\eta \end{pmatrix} \leq A_n + \eta A_n^2 \quad \text{for all } \eta > 0,$$

where

$$A_n := n^2 \begin{pmatrix} I_d & -I_d \\ -I_d & I_d \end{pmatrix} + Q.$$

Theorem A.1. *Let Assumption E hold true. Let w_1 be an u.s.c. viscosity subsolution on $\overline{\mathcal{O}}$ and let w_2 be an l.s.c. viscosity supersolution on \mathcal{O} of (A.1). If w_1 and w_2 have polynomial growth on $\overline{\mathcal{O}}$ and if w_2 is of class $\mathcal{R}(\mathcal{O})$, then*

$$w_2 \geq w_1 \text{ on } \{T\} \times \overline{\mathcal{O}} \quad \text{implies} \quad w_2 \geq w_1 \text{ on } [0, T] \times \overline{\mathcal{O}}.$$

Remark A.2. Our result also applies to the equation

$$-\partial_t \varphi + \mathcal{H}(\cdot, D\varphi, D^2\varphi) = 0, \quad (\text{A.2})$$

provided that \mathcal{H} is homogeneous of degree one with respect to its second and third argument, as it is the case for the Hamilton-Jacobi-Bellman operators in the body of this paper. Indeed, w_1 is then a subsolution of (A.2) if and only if $(t, x) \mapsto e^{\rho t} w_1(t, x)$ is a subsolution of (A.1), and similarly for the supersolution. Further extensions could also be considered but are beyond the scope of this paper.

Proof of Theorem A.1. Assume that $w_2 \geq w_1$ on $\{T\} \times \overline{\mathcal{O}}$. Let $p \geq 1$ and $C > 0$ be such that $w_1(t, x) - w_2(t, x) \leq C(1 + |x|^p)$ for all $(t, x) \in [0, T] \times \overline{\mathcal{O}}$. Assume for contradiction that $\sup(w_1 - w_2) > 0$, then we can find $\iota > 0$ and $(t_0, x_0) \in [0, T] \times \overline{\mathcal{O}}$ such that

$$\zeta := (w_1 - 2\phi - w_2)(t_0, x_0) = \max_{[0, T] \times \overline{\mathcal{O}}} (w_1 - 2\phi - w_2) > 0, \quad (\text{A.3})$$

where

$$\phi(t, x) := \iota e^{-\kappa t} (1 + |x|^{2p}).$$

Here $\kappa > 0$ is a fixed constant which is large enough to ensure that

$$\begin{aligned} m(t, x) := \\ -\rho\phi(t, x) + \partial_t\phi(t, x) + \alpha((1 + |x|)|D\phi(t, x)| + (1 + |x|^2)|D^2\phi(t, x)|) \end{aligned} \quad (\text{A.4})$$

is nonpositive on $[0, T] \times \mathbb{R}^d$, where α is the constant from Assumption E.

Note that by the assumption that $w_2 \geq w_1$ on $\{T\} \times \overline{\mathcal{O}}$, we must have $(t_0, x_0) \in [0, T] \times \overline{\mathcal{O}}$.

Case 1: $x_0 \in \partial\mathcal{O}$. For all $n \geq 1$, there exist $(t^n, x^n, s^n, y^n) \in ([0, T] \times \overline{\mathcal{O}})^2$ satisfying

$$\Phi^n(t^n, x^n, s^n, y^n) = \max_{([0, T] \times \overline{\mathcal{O}})^2} \Phi^n, \quad (\text{A.5})$$

where

$$\Phi^n(t, x, s, y) := w_1(t, x) - w_2(s, y) - \Theta^n(t, x, s, y)$$

and

$$\begin{aligned} \Theta^n(t, x, s, y) := & \frac{1}{2} n^2 (|t + \lambda(n^{-1}) - s|^2 + \varepsilon |x + \ell(n^{-1}) - y|^2) \\ & + |t - t_0|^2 + |x - x_0|^4 + \phi(t, x) + \phi(s, y), \end{aligned}$$

with ℓ and λ given for x_0 as in the statement of the Definition 4.10, and $\varepsilon > 0$. Note that (A.5), the assumption that w_2 satisfies (4.22), and (A.3) imply that

$$\begin{aligned} \Phi^n(t^n, x^n, s^n, y^n) & \geq \Phi^n(t_0, x_0, t_0 + \lambda(n^{-1}), x_0 + \ell(n^{-1})) \\ & = (w_1 - 2\phi - w_2)(t_0, x_0) + o(1) \\ & = \zeta + o(1) \quad \text{as } n \rightarrow \infty. \end{aligned} \quad (\text{A.6})$$

Recalling the growth condition on w_1, w_2 and the definition of ϕ , it follows that, after passing to a subsequence, (t^n, x^n, s^n, y^n) converges to some $(t^\infty, x^\infty, t^\infty, x^\infty) \in ([0, T] \times \overline{\mathcal{O}})^2$. We then have

$$\begin{aligned} \zeta & = (w_1 - 2\phi - w_2)(t_0, x_0) \\ & = \max_{[0, T] \times \overline{\mathcal{O}}} (w_1 - 2\phi - w_2) \\ & \geq (w_1 - 2\phi - w_2)(t^\infty, x^\infty) - |t^\infty - t_0|^2 - |x^\infty - x_0|^4 \\ & \quad - \limsup_{n \rightarrow \infty} \frac{1}{2} n^2 (|t^n + \lambda(n^{-1}) - s^n|^2 + \varepsilon |x^n + \ell(n^{-1}) - y^n|^2) \\ & \geq \liminf_{n \rightarrow \infty} \Phi^n(t^n, x^n, s^n, y^n) \\ & \geq \zeta, \end{aligned}$$

where (A.6) was used in the last step. After passing to a subsequence, we deduce that

$$(t^n, x^n, s^n, y^n) \rightarrow (t_0, x_0, t_0, x_0), \quad (\text{A.7})$$

$$w_1(t^n, x^n) - w_2(s^n, y^n) \rightarrow (w_1 - w_2)(t_0, x_0), \quad (\text{A.8})$$

$$s^n = t^n + \lambda(n^{-1}) + o(n^{-1}), \quad y^n = x^n + \ell(n^{-1}) + o(n^{-1}). \quad (\text{A.9})$$

Since $(t_0, x_0) \in [0, T] \times \partial\mathcal{O}$, it follows from (4.20), (4.21) and (A.9) that $(s^n, y^n) \in [0, T] \times \mathcal{O}$ for n large enough.

Let $\overline{\mathcal{P}}_{\mathcal{O}}^{2,+} w_1$ and $\overline{\mathcal{P}}_{\mathcal{O}}^{2,-} w_2$ be the “closed” parabolic super- and subjets as defined in [5, Section 8]. From the Crandall-Ishii lemma [5, Theorem 8.3] we obtain, for each $\eta > 0$, elements

$$(a^n, p^n, X_\eta^n) \in \overline{\mathcal{P}}_{\mathcal{O}}^{2,+} w_1(t^n, x^n) \quad \text{and} \quad (b^n, q^n, Y_\eta^n) \in \overline{\mathcal{P}}_{\mathcal{O}}^{2,-} w_2(s^n, y^n)$$

such that

$$\begin{aligned} a^n &= \partial_t \Theta^n(t^n, x^n, s^n, y^n), & b^n &= -\partial_s \Theta^n(t^n, x^n, s^n, y^n), \\ p^n &= D_x \Theta^n(t^n, x^n, s^n, y^n), & q^n &= -D_y \Theta^n(t^n, x^n, s^n, y^n), \end{aligned}$$

$$\begin{pmatrix} X_\eta^n & 0 \\ 0 & -Y_\eta^n \end{pmatrix} \leq A_n + \eta A_n^2,$$

where $A_n = D^2 \Theta^n(t^n, x^n, s^n, y^n)$; i.e.,

$$A_n = \varepsilon n^2 \begin{pmatrix} I_d & -I_d \\ -I_d & I_d \end{pmatrix} + \begin{pmatrix} D^2 \phi(t^n, x^n) + O(|x^n - x_0|^2) & 0 \\ 0 & D^2 \phi(s^n, y^n) \end{pmatrix}.$$

In view of the sub- and supersolution properties of w_1 and w_2 , the fact that $(s^n, y^n) \in [0, T) \times \mathcal{O}$ for n large, and Assumption E, we deduce that

$$\begin{aligned} \Delta_n &:= \rho(w_1(t^n, x^n) - w_2(s^n, y^n)) \\ &\leq 2(t^n - t_0) + \partial_t \phi(t^n, x^n) + \partial_s \phi(s^n, y^n) \\ &\quad + \alpha |x^n - y^n| (1 + |q^n| + \varepsilon n^2 |x^n - y^n|) \\ &\quad + \alpha ((1 + |x^n|) |q^n - p^n| + (1 + |x^n|^2) |Q^n|), \end{aligned}$$

where

$$Q^n := \begin{pmatrix} D^2 \phi(t^n, x^n) + O(|x^n - x_0|^2) & 0 \\ 0 & D^2 \phi(s^n, y^n) \end{pmatrix}.$$

By the definitions of p^n and q^n , it follows that

$$\begin{aligned} \Delta_n &\leq 2(t^n - t_0) + \partial_t \phi(t^n, x^n) + \partial_s \phi(s^n, y^n) \\ &\quad + \alpha |x^n - y^n| (1 + |D\phi(s^n, y^n)| + \varepsilon n^2 |x^n + \ell(n^{-1}) - y^n| + \varepsilon n^2 |x^n - y^n|) \\ &\quad + \alpha (1 + |x^n|) (4|x^n - x_0|^3 + |D\phi(t^n, x^n)| + |D\phi(s^n, y^n)|) \\ &\quad + \alpha (1 + |x^n|^2) (|D^2 \phi(t^n, x^n)| + |D^2 \phi(s^n, y^n)| + O(|x^n - x_0|^2)). \end{aligned}$$

Recalling (A.7)–(A.9), letting $n \rightarrow \infty$ leads to

$$\begin{aligned} \rho(w_1 - w_2)(t_0, x_0) &\leq 2\partial_t \phi(t_0, x_0) + \alpha \varepsilon \left(\liminf_{n \rightarrow \infty} n \ell(n^{-1}) \right)^2 \\ &\quad + 2\alpha ((1 + |x_0|) |D\phi(t_0, x_0)| + (1 + |x_0|^2) |D^2 \phi(t_0, x_0)|), \end{aligned}$$

which, by (4.19) and the definition of m in (A.4), implies

$$\rho(w_1 - 2\phi - w_2)(t_0, x_0) \leq 2m(t_0, x_0)$$

after letting $\varepsilon \rightarrow 0$. Since $\kappa > 0$ has been chosen so that $m \leq 0$ on $[0, T] \times \mathbb{R}^d$, this contradicts (A.3).

Case 2: $x_0 \in \mathcal{O}$. This case is handled similarly by using

$$\Theta^n(t, x, s, y) := \frac{1}{2}n^2(|t-s|^2 + |x-y|^2) + |t-t_0|^2 + |x-x_0|^4 + \phi(t, x) + \phi(s, y).$$

After taking a subsequence, the corresponding sequence of maximum points $(t^n, x^n, s^n, y^n)_{n \geq 1}$ again converges to (t_0, x_0, t_0, x_0) , so that $x^n, y^n \in \mathcal{O}$ for n large enough. The rest of the proof follows the same arguments as in Case 1. \square

Lemma A.3. *Under Assumption D', the operator \bar{H} defined in (4.17) satisfies Assumption E.*

Proof. The argument is standard. We first observe that

$$\begin{aligned} -\bar{L}^u(y, q, Y^\eta) + \bar{L}^u(x, p, X^\eta) &= (\mu(x, u) - \mu(y, u))^\top q + \mu(x, u)^\top (p - q) \\ &\quad + \frac{1}{2} \sum_{1 \leq i \leq d} \Sigma^i(x, y, u)^\top \Xi^\eta \Sigma^i(x, y, u), \end{aligned}$$

where

$$\Sigma(x, y, u) := \begin{pmatrix} \sigma(x, u) \\ \sigma(y, u) \end{pmatrix} \quad \text{and} \quad \Xi^\eta := \begin{pmatrix} X^\eta & 0 \\ 0 & -Y^\eta \end{pmatrix}$$

and Σ^i denotes the i th column of Σ . Since μ is Lipschitz continuous and has uniformly linear growth by Assumption D', we have

$$(\mu(x, u) - \mu(y, u))^\top q + \mu(x, u)^\top (p - q) \leq \alpha(|x - y||q| + (1 + |x|)|p - q|)$$

for some constant $\alpha > 0$. Recall from the statement of Assumption E the condition that

$$\Xi^\eta \leq A_n + \eta A_n^2, \quad \text{where} \quad A_n := n^2 \begin{pmatrix} I_d & -I_d \\ -I_d & I_d \end{pmatrix} + Q.$$

Since σ is Lipschitz continuous and satisfies Assumption D', it follows that

$$\begin{aligned} \Sigma^i(x, y, u)^\top \Xi^\eta \Sigma^i(x, y, u) &\leq n^2 |\sigma^i(x, u) - \sigma^i(y, u)|^2 \\ &\quad + \alpha(1 + |x|^2 + |y|^2)|Q| + O(\eta), \end{aligned}$$

possibly after enlarging α , where σ^i denotes the i th column of σ . We conclude by using the Lipschitz continuity of σ and the condition that $|x - y| \leq 1$. \square

Remark A.4. We conclude with some comments on Theorem A.1 and Proposition 4.12, and related results in the literature.

The main issue in proving comparison with state constraints is to avoid boundary points; i.e., that y_n (see the proof of Theorem A.1) ends up on the boundary. One classical way to ensure this, is to use a perturbation of $|x - y|^2$ in a suitable inward direction, like the function ℓ above. Moreover, this requires the supersolution to be continuous at the boundary points, along the direction of perturbation; cf. (4.22).

In [11], the inner normal $n(x)$ at the boundary point $x \in \partial\mathcal{O}$ is used as an inward direction. In the proof, one is only close to x (cf. (A.9)); therefore, the comparison result [11, Theorem 2.2] requires the existence of a truncated cone around $n(x)$ which stays inside the domain, and the continuity of the subsolution along the directions that it generates. Our condition (4.20) is less restrictive than the corresponding requirement in [11]: we only need the continuity along the curve $\varepsilon \mapsto \ell(\varepsilon)$, cf. (4.22), rather than all lines in a neighborhood. The function λ appears because we consider parabolic equations, whereas [11] focuses on the elliptic case.

In Proposition 4.12 we give conditions (certainly not the most general possible) ensuring that the value function is of class $\mathcal{R}(\mathcal{O})$. They should be compared to [11, Condition (A3)], which is used to verify the continuity assumption of [11, Theorem 2.2]. Our conditions are stronger in the sense that they are imposed around the boundary and not only at the boundary; on the other hand, we require C^1 -regularity of the boundary whereas [11] requires C^3 .

In [10], a slightly different technique is used, based on ideas from [9]. First, it is assumed that at each boundary point x , there exists a fixed control which kills the volatility at the neighboring boundary points and keeps a truncated cone around the drift in the domain. This is similar to our (4.23), except that our control is not fixed; on the other hand, we assume that it kills the volatility in a neighborhood of x . Thus, the conditions in [10] are not directly comparable to ours; e.g., if \mathcal{O} is the unit disk in \mathbb{R}^2 , $U = [-1, 1]$, $\mu(x, u) = -x$ and $\sigma(x, u) = |x^2 - u|I_2$, then (4.23) is satisfied (with δ given by the Euclidean distance near the boundary and $\check{u}(x) = x^2$), while [10, Condition (2.1)] is not. Second, in [10], the state constraint problem is transformed so as to introduce a Neumann-type boundary condition and construct a suitable test function which, as a function of its first component, turns out to be a uniformly strict supersolution of the Neumann boundary condition. The construction in [10] heavily relies on the assumption that the coefficients are bounded; cf. the beginning of [10, Section 3] and the proof of [10, Theorem 3.1].

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